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A new flocking model through body attitude coordination

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We present a new model for multi-agent dynamics where each agent is described by its position and body attitude: agents travel at a constant speed in a given direction and their body can rotate around it adopting different configurations. In this manner, the body attitude is described by three orthonormal axes giving an element in $SO(3)$ (rotation matrix). Agents try to coordinate their body attitudes with the ones of their neighbours. In this paper, we give the individual-based model (particle model) for this dynamics and derive its corresponding kinetic and macroscopic equations.

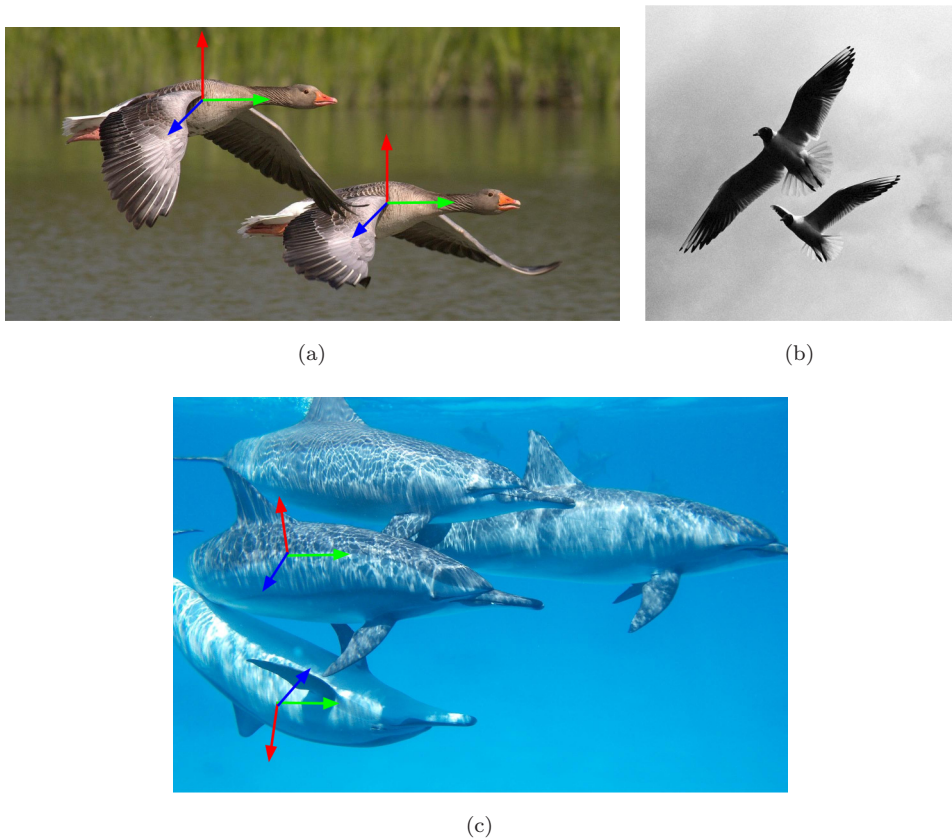
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1. Introduction

In this paper, we model collective motion where individuals or agents are described by their position and body attitude. The body attitude is given by three

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Fig. 1. (Color online) Examples of body attitude coordination/dis-coordination and the use of the rotation matrix representation. (a) Birds with coordinated body attitude. Three orthonormal axes describe the body attitude: the green arrow indicates the direction of movement; the blue and red arrows indicate the position of the body with respect to this direction. (b) Birds with no coordinated body attitude. (c) Dolphins moving in the same direction but with different body attitude. In this example one can see that the body attitude coordination model gives more information than the Vicsek model (which only describes the direction of movement).

orthonormal axes; one of the axes describes the direction in which the agent moves at a constant speed; the other two axes indicate the relative position of the body with respect to this direction. Agents try to coordinate their body attitude with those of near neighbours (see Fig. 1). Here we present an individual-based model (particle model) for body attitude coordination and derive the corresponding macroscopic equations from the associated mean-field equation, which we refer to as the Self-Organised Hydrodynamics for body attitude coordination (SOHB), by reference to the Self-Organised Hydrodynamics (SOH) derived from the Vicsek dynamics (see Ref. 24 and discussion below).

There exists already a variety of models for collective behaviour depending on the type of interaction between agents. However, to the best of our knowledge, this is the first model that takes into account interactions based on body attitude coordination. This has applications in the study of collective motion of animals such as birds and fish and it is a stepping stone to model more complex agents formed by articulated bodies (corpora).^{13,14} In this section, we present related results in the literature and the structure of the document.

There exists a vast literature on collective behaviour. In particular, here we deal with the case of self-propelled particles which is ubiquitous in nature. It includes, among others, fish schools, flocks of birds, herds^{8,9,41}; bacteria^{5,47}; human walking behaviour.³⁶ The interest in studying these systems is to gain understanding on the emergent properties: local interactions give rise to large scale structures in the form of patterns and phase transitions (see the review in Ref. 46). These large scale structures take place when the number of individuals or agents is very large. In this case a statistical description of the system is more pertinent than an agent-based one. With this in mind mean-field limits are devised when the number of agents tend to infinity. From them macroscopic equations can be obtained using hydrodynamic limit techniques (as we explain below).

The results presented here are inspired from those in Ref. 24. There the authors consider the Vicsek model which is a particular type of model for self-propelled particles.^{1,15,34,45} The Vicsek model describes collective motion where agents travel at a constant speed in a given direction. At each time step the direction of movement is updated to the averaged one of the neighbouring agents, with some noise. The position is updated considering the distance travelled during that time step.

Our results here are inspired by the SOH model (the continuum version of the Vicsek model), where we have substituted velocity alignment by body attitude coordination. Other refinements and adaptations of the Vicsek model (at the particle level) or the SOH model (at the continuum level) have been proposed in the literature, we just mention the following ones as examples: in Ref. 10 an individual-based model is proposed to better describe collective motion of turning birds; in Ref. 25 agents are considered to have the shape of discs and volume exclusion is included in the dynamics; a description of nematic alignment in rods is done in Ref. 23.

In Ref. 24 the authors investigate the mean-field limit and macroscopic limit of the Vicsek model. The mean-field limit gives a kinetic equation that takes the form of a Fokker–Planck equation referred to as the mean-field limit Vicsek model.

To obtain the macroscopic equations (the SOH model), the authors in Ref. 24 use the well-known tools of hydrodynamic limits, first developed in the framework of the Boltzmann equation for rarefied gases.^{11,16,42} Since its first appearance, hydrodynamics limits have been used in other different contexts, including traffic flow modelling^{3,35} and supply chain research.^{2,26} However, in Ref. 24 a new concept is introduced: the Generalised Collision Invariant (GCI). Typically to obtain the macroscopic equations we require as many conserved quantities in the kinetic

equation as the dimension of the equilibria (see again Ref. 46). In the mean-field limit Vicsek model this requirement is not fulfilled and the GCI is used to sort out this problem. For other cases where the GCI concept has been used see Refs. 17, 18, 21, 22, 27 and 29.

After this introduction and the following discussion about the main result, the next part of the paper is dedicated to the modelling. In Sec. 3.1 we explain the derivation of the individual-based model for body coordination dynamics: given N agents labelled by $k = 1, \dots, N$ the positions and body attitudes $(\mathbf{X}_k, A_k) \in \mathbb{R}^3 \times \text{SO}(3)$ over time are given by the Stochastic Differential Equations (SDEs) (3.10)–(3.11). In Sec. 3.2 we give the corresponding (formal) mean-field limit (Proposition 3.2) for the evolution of the empirical measure when the number of agents $N \rightarrow \infty$.

The last part concerns the derivation of the macroscopic equations (Theorem 4.1) for the total density of the particles $\rho = \rho(t, x)$ and the matrix of the mean body attitude $\Lambda = \Lambda(t, x)$. To obtain these equations we first study the rescaled mean-field equation (Eq. (4.1)), which is, at leading order, a Fokker–Planck equation. We determine its equilibria, which are given by a von Mises distribution on $\text{SO}(3)$ (Eq. (4.4)). Finally in Sec. 4.4 we obtain the GCIs (Proposition 4.6), which are the main tool to be able to derive the macroscopic equations in Sec. 4.5.

2. Discussion of the Main Result: The Self-Organised Hydrodynamics for Body Attitude Coordination

The main result of this paper is Theorem 4.1 which gives the following macroscopic equations for the density of agents $\rho = \rho(t, x)$ and the matrix of the mean body attitude $\Lambda = \Lambda(t, x) \in \text{SO}(3)$ (i.e. SOHB):

$$\partial_t \rho + c_1 \nabla_x \cdot (\rho \Lambda \mathbf{e}_1) = 0, \quad (2.1)$$

$$\begin{aligned} \rho(\partial_t \Lambda + c_2((\Lambda \mathbf{e}_1) \cdot \nabla_x) \Lambda) + [(\Lambda \mathbf{e}_1) \\ \times (c_3 \nabla_x \rho + c_4 \rho \mathbf{r}_x(\Lambda)) + c_4 \rho \delta_x(\Lambda) \Lambda \mathbf{e}_1]_{\times} \Lambda = 0. \end{aligned} \quad (2.2)$$

In the equations above c_1 , c_2 , c_3 and c_4 are explicit constants which depend on the parameters of the model (namely the rate of coordination and the level of noise). The expressions of the constants c_2 , c_3 and c_4 depend on the GCI mentioned in the Introduction (and computed thanks to Proposition 4.6). The constant c_1 is obtained as a “consistency” relation (Lemma 4.4). In (2.2), the operation $[\cdot]_{\times}$ transforms a vector \mathbf{v} in an antisymmetric matrix such that $[\mathbf{v}]_{\times} \mathbf{u} = \mathbf{v} \times \mathbf{u}$ for any vector \mathbf{u} (see (3.2) for the exact definition). The scalar $\delta_x(\Lambda)$ and the vector $\mathbf{r}_x(\Lambda)$ are first-order differential operators intrinsic to the dynamics: if $\Lambda(x) = \exp([\mathbf{b}(x)]_{\times}) \Lambda(x_0)$ with \mathbf{b} smooth around x_0 and $\mathbf{b}(x_0) = 0$, then

$$\delta_x(\Lambda)(x_0) = \nabla_x \cdot \mathbf{b}(x)|_{x=x_0} \quad \text{and} \quad \mathbf{r}_x(\Lambda)(x_0) = \nabla_x \times \mathbf{b}(x)|_{x=x_0},$$

where $\nabla_x \times$ is the curl operator. These operators are well defined as long as Λ is smooth: as we will see in the next section, we can always express a rotation matrix

as $\exp([\mathbf{b}]_{\times})$ for some vector $\mathbf{b} \in \mathbb{R}^3$, and this function $\mathbf{b} \mapsto \exp([\mathbf{b}]_{\times})$ is a local diffeomorphism between a neighbourhood of $0 \in \mathbb{R}^3$ and the identity of $\text{SO}(3)$. This gives a unique smooth representation of \mathbf{b} in the neighbourhood of 0 when x is in the neighbourhood of x_0 since then $\Lambda(x)\Lambda(x_0)^{-1}$ is in the neighbourhood of Id .

We express Eq. (2.2) in terms of the basis vectors $\{\Omega = \Lambda \mathbf{e}_1, \mathbf{u} = \Lambda \mathbf{e}_2, \mathbf{v} = \Lambda \mathbf{e}_3\}$ and we write $\Lambda = \Lambda(\Omega, \mathbf{u}, \mathbf{v})$. System (2.1)–(2.2) can be expressed as an evolution system for ρ and the basis $\{\Omega, \mathbf{u}, \mathbf{v}\}$ as follows:

$$\partial_t \rho + c_1 \nabla_x \cdot (\rho \Omega) = 0, \quad (2.3)$$

$$\rho D_t \Omega + P_{\Omega^\perp} (c_3 \nabla_x \rho + c_4 \rho \mathbf{r}) = 0, \quad (2.4)$$

$$\rho D_t \mathbf{u} - \mathbf{u} \cdot (c_3 \nabla_x \rho + c_4 \rho \mathbf{r}) \Omega + c_4 \rho \delta \mathbf{v} = 0, \quad (2.5)$$

$$\rho D_t \mathbf{v} - \mathbf{v} \cdot (c_3 \nabla_x \rho + c_4 \rho \mathbf{r}) \Omega - c_4 \rho \delta \mathbf{u} = 0, \quad (2.6)$$

where $D_t := \partial_t + c_2(\Omega \cdot \nabla_x)$, $\delta = \delta_x(\Lambda(\Omega, \mathbf{u}, \mathbf{v}))$ and $\mathbf{r} = \mathbf{r}_x(\Lambda(\Omega, \mathbf{u}, \mathbf{v}))$. The operator P_{Ω^\perp} denotes the projection on the orthogonal of Ω . We easily see here that these equations preserve the constraints $|\Omega| = |\mathbf{u}| = |\mathbf{v}| = 1$ and $\Omega \cdot \mathbf{u} = \Omega \cdot \mathbf{v} = \mathbf{u} \cdot \mathbf{v} = 0$. The expressions of δ and \mathbf{r} are:

$$\delta = [(\Omega \cdot \nabla_x) \mathbf{u}] \cdot \mathbf{v} + [(\mathbf{u} \cdot \nabla_x) \mathbf{v}] \cdot \Omega + [(\mathbf{v} \cdot \nabla_x) \Omega] \cdot \mathbf{u},$$

$$\mathbf{r} = (\nabla_x \cdot \Omega) \Omega + (\nabla_x \cdot \mathbf{u}) \mathbf{u} + (\nabla_x \cdot \mathbf{v}) \mathbf{v}.$$

Equation (2.1) (or equivalently Eq. (2.3)) is the continuity equation for ρ and ensures mass conservation. The convection velocity is given by $c_1 \Lambda \mathbf{e}_1 = c_1 \Omega$ and Ω gives the direction of motion. Equation (2.2) (or equivalently Eqs. (2.4)–(2.6)) gives the evolution of Λ . We remark that every term in Eq. (2.2) belongs to the tangent space at Λ in $\text{SO}(3)$; this is true for the first term since $(\partial_t + c_2(\Lambda \mathbf{e}_1) \cdot \nabla_x)$ is a differential operator and it also holds for the second term because it is the product of an antisymmetric matrix with Λ (see Proposition A.2). Alternately, this means that $(\Omega(t), \mathbf{u}(t), \mathbf{v}(t))$ is a direct orthonormal basis as soon as $(\Omega(0), \mathbf{u}(0), \mathbf{v}(0))$.

The term corresponding to c_3 in (2.2) gives the influence of $\nabla_x \rho$ (pressure gradient) on the body attitude Λ . It has the effect of rotating the body around the vector directed by $(\Lambda \mathbf{e}_1) \times \nabla_x \rho$ at an angular speed given by $\frac{c_3}{\rho} \|(\Lambda \mathbf{e}_1) \times \nabla_x \rho\|$, so as to align Ω with $-\nabla_x \rho$. Indeed the solution of the differential equation $\frac{d\Lambda}{dt} + \gamma[\mathbf{n}]_{\times} \Lambda = 0$, when \mathbf{n} is a constant unit vector and γ a constant scalar, is given by $\Lambda(t) = \exp(-\gamma t[\mathbf{n}]_{\times}) \Lambda_0$, and $\exp(-\gamma t[\mathbf{n}]_{\times})$ is the rotation of axis \mathbf{n} and angle $-\gamma t$ (see (3.4), it is called Rodrigues' formula). Since we will see that c_3 is positive the influence of this term consists of relaxing the direction of displacement $\Lambda \mathbf{e}_1$ towards $\nabla_x \rho$. Alternately, we can see from (2.4) that Ω turns in the opposite direction to $\nabla_x \rho$, showing that the $\nabla_x \rho$ term has the same effect as a pressure gradient in classical hydrodynamics. We note that the pressure gradient has also the effect of rotating the whole body frame (see influence of $\nabla_x \rho$ on \mathbf{u} and \mathbf{v}) just to keep the frame orthonormal. Similarly to what happens with the $\nabla_x \rho$ term in Eq. (2.2), the term containing $c_4 \rho \mathbf{r}$ in Eq. (2.4) has the effect of relaxing the direction of

displacement Ω towards $-\mathbf{r}$ (we will indeed see that c_4 is positive). Finally, the last terms of Eqs. (2.5)–(2.6) have the effect of rotating the vectors \mathbf{u} and \mathbf{v} around Ω along the flow driven by D_t at angular speed $c_4\delta$.

If we forget the term proportional to \mathbf{r} in (2.4), system (2.3)–(2.4) is decoupled from (2.5)–(2.6), and is an autonomous system for ρ and Ω , which coincides with the SOH model. The SOH model provides the fluid description of a particle system obeying the Vicsek dynamics.²⁴ As already discussed in Ref. 24, the SOH model bears analogies with the compressible Euler equations, where (2.3) is obviously the mass conservation equation and (2.4) is akin to the momentum conservation equation, where momentum transport $\rho D_t \Omega$ is balanced by a pressure force $-P_{\Omega^\perp} \nabla_x \rho$. There are however major differences. The first one is the presence of the projection operation P_{Ω^\perp} which is there to preserve the constraint $|\Omega| = 1$. Indeed, while the velocity in the Euler equations is an arbitrary vector, the quantity Ω in the SOH model is a velocity orientation and is normalised to 1. The second one is that the convection speed c_2 in the convection operator D_t is *a priori* different from the mass convection speed c_1 appearing in the continuity equation. This difference is a signature of the lack of Galilean invariance of the system, which is a common feature of all dry active matter models.

The major novelty of the present model, which can be referred to as the SOHB is that the transport of the direction of motion Ω involves the influence of another quantity specific to the body orientation dynamics, namely the vector \mathbf{r} . The overall dynamics tends to align the velocity orientation Ω , not opposite to the density gradient $\nabla_x \rho$ but opposite to a composite vector $(c_3 \nabla_x \rho + c_4 \rho \mathbf{r})$. The vector \mathbf{r} is the rotational of a vector \mathbf{b} locally attached to the frame (namely the unit vector of the local rotation axis multiplied by the local angle of rotation around this axis). This vector gives rise to an effective pressure force which adds up to the usual pressure gradient. It would be interesting to design specific solutions where this effective pressure force has a demonstrable effect on the velocity orientation dynamics.

In addition to this effective force, spatial inhomogeneities of the body attitude also have the effect of inducing a proper rotation of the frame about the direction of motion. This proper rotation is also driven by spatial inhomogeneities of the vector \mathbf{b} introduced above, but are now proportional to its divergence.

3. Modelling: The Individual-Based Model and Its Mean-Field Limit

The body attitude is given by a rotation matrix. Therefore, we work on the Riemannian manifold $\text{SO}(3)$ (Special Orthogonal Group), which is formed by the subset of matrices A such that $AA^T = \text{Id}$ and $\det(A) = 1$, where Id stands for the identity matrix.

In this document \mathcal{M} indicates the set of square matrices of dimension 3; \mathcal{A} is the set of antisymmetric matrices of dimension 3; \mathcal{S} is the set of symmetric matrices of

dimension 3. Typically we will denote by A, Λ matrices in $\text{SO}(3)$ and by P matrices in \mathcal{A} . Bold symbols $\mathbf{n}, \mathbf{v}, \mathbf{e}_1$ indicate vectors.

We will often use the so-called Euler axis-angle parameters to represent an element in $\text{SO}(3)$: to $A \in \text{SO}(3)$ there is associated an angle $\theta \in [0, \pi]$ and a vector $\mathbf{n} \in S^2$ so that $A = A(\theta, \mathbf{n})$ corresponds to the anticlockwise rotation of angle θ around the vector \mathbf{n} . It is easy to see that

$$\text{tr}(A) = 1 + 2 \cos \theta \quad (3.1)$$

(for instance expressing A in an orthonormal basis with \mathbf{n}), so the angle θ is uniquely defined as $\arccos(\frac{1}{2}(\text{tr}(A) - 1))$. Notice that \mathbf{n} is uniquely defined whenever $\theta \in (0, \pi)$ (if $\theta = 0$ then \mathbf{n} can be any vector in S^2 and if $\theta = \pi$ then the direction of \mathbf{n} is uniquely defined but not its orientation). For a given vector \mathbf{u} , we introduce the antisymmetric matrix $[\mathbf{u}]_{\times}$, where $[\cdot]_{\times}$ is the linear operator from \mathbb{R}^3 to \mathcal{A} given by

$$[\mathbf{u}]_{\times} := \begin{bmatrix} 0 & -u_3 & u_2 \\ u_3 & 0 & -u_1 \\ -u_2 & u_1 & 0 \end{bmatrix}, \quad (3.2)$$

so that for any vectors $\mathbf{u}, \mathbf{v} \in \mathbb{R}^3$, we have $[\mathbf{u}]_{\times} \mathbf{v} = \mathbf{u} \times \mathbf{v}$. In this framework, we have the following representation for $A \in \text{SO}(3)$ (called Rodrigues' formula):

$$A = A(\theta, \mathbf{n}) = \text{Id} + \sin \theta [\mathbf{n}]_{\times} + (1 - \cos \theta) [\mathbf{n}]_{\times}^2 \quad (3.3)$$

$$= \exp(\theta [\mathbf{n}]_{\times}). \quad (3.4)$$

We also have $\mathbf{n} \times (\mathbf{n} \times \mathbf{v}) = (\mathbf{n} \cdot \mathbf{v})\mathbf{n} - (\mathbf{n} \cdot \mathbf{n})\mathbf{v}$, therefore when \mathbf{n} is a unit vector, we have

$$[\mathbf{n}]_{\times}^2 = \mathbf{n} \otimes \mathbf{n} - \text{Id}, \quad (3.5)$$

where the tensor product $\mathbf{a} \otimes \mathbf{b}$ is the matrix defined by $(\mathbf{a} \otimes \mathbf{b})\mathbf{u} = (\mathbf{u} \cdot \mathbf{b})\mathbf{a}$ for any $\mathbf{u} \in \mathbb{R}^3$. Finally, $\text{SO}(3)$ has a natural Riemannian metric (see Ref. 38) induced by the following inner product in the set of square matrices of dimension 3:

$$A \cdot B = \frac{1}{2} \text{tr}(A^T B) = \frac{1}{2} \sum_{i,j} A_{ij} B_{ij}. \quad (3.6)$$

This normalisation gives us that for any vectors $\mathbf{u}, \mathbf{v} \in \mathbb{R}^3$, we have that

$$[\mathbf{u}]_{\times} \cdot [\mathbf{v}]_{\times} = (\mathbf{u} \cdot \mathbf{v}). \quad (3.7)$$

Moreover, the geodesic distance on $\text{SO}(3)$ between Id and a rotation of angle $\theta \in [0, \pi]$ is exactly given by θ (the geodesic between Id and A is exactly $t \in [0, \theta] \mapsto \exp(t[\mathbf{n}]_{\times})$). See Appendix A for some properties of $\text{SO}(3)$ used throughout this work.

Seeing $\text{SO}(3)$ as a Riemannian manifold, we will use the following notations: T_A is the tangent space in $\text{SO}(3)$ at $A \in \text{SO}(3)$; P_{T_A} denotes the orthogonal projection onto T_A ; the operators $\nabla_A, \nabla_A \cdot$ are the gradient and divergence in $\text{SO}(3)$, respectively. These operators are computed in Sec. 4.2 in the Euler axis-angle coordinates.

3.1. The individual-based model

Consider N agents labelled by $k = 1, \dots, N$ with positions $\mathbf{X}_k(t) \in \mathbb{R}^3$ and associated matrices (body attitudes) $A_k(t) \in \text{SO}(3)$. For each k , the three unit vectors representing the frame correspond to the vectors of the matrix $A_k(t)$ when written as a matrix in the canonical basis $(\mathbf{e}_1, \mathbf{e}_2, \mathbf{e}_3)$ of \mathbb{R}^3 . In particular, the direction of displacement of the agent is given by its first vector $A_k(t)\mathbf{e}_1$.

Evolution of the positions. Agents move in the direction of the first axis with constant speed v_0 :

$$\frac{d\mathbf{X}_k(t)}{dt} = v_0 A_k(t)\mathbf{e}_1.$$

Evolution of the body attitude matrix. Agents try to coordinate their body attitude with those of their neighbours. So we are facing two different problems from a modelling viewpoint, namely to define the target body attitude, and to define the way agents relax their own attitude towards this “average” attitude.

As for the Vicsek model,²⁴ we consider a kernel of influence $K = K(x) \geq 0$ and define the matrix

$$\mathbb{M}_k(t) := \frac{1}{N} \sum_{i=1}^N K(|\mathbf{X}_i(t) - \mathbf{X}_k(t)|) A_i(t). \quad (3.8)$$

This matrix corresponds to the averaged body attitude of the agents inside the zone of influence corresponding to agent k . Now $\mathbb{M}_k(t) \notin \text{SO}(3)$, so we need to orthogonalise and remove the dilations, in order to construct a target attitude in $\text{SO}(3)$. We will see that the polar decomposition of $\mathbb{M}_k(t)$ is a good choice in the sense that it minimises a weighted sum of the squared distances to the attitudes of the neighbours. We also refer to Ref. 40 for some complements on averaging in $\text{SO}(3)$.

We give next the definition of polar decomposition.

Lemma 3.1. (Polar decomposition of a square matrix³³) *Given a matrix $M \in \mathcal{M}$, if $\det(M) \neq 0$ then there exists a unique orthogonal matrix A (given by $A = M(\sqrt{M^T M})^{-1}$) and a unique symmetric positive definite matrix S such that $M = AS$.*

Proposition 3.1. *Suppose that the matrix $\mathbb{M}_k(t)$ has positive determinant. Then the following assertions are equivalent:*

- (i) *The matrix A minimises the quantity $\frac{1}{N} \sum_{i=1}^N K(|\mathbf{X}_i(t) - \mathbf{X}_k(t)|) \|A_i(t) - A\|^2$ among the elements of $\text{SO}(3)$.*
- (ii) *The matrix A is the element of $\text{SO}(3)$ which maximises the quantity $A \cdot \mathbb{M}_k(t)$.*
- (iii) *The matrix A is the polar decomposition of $\mathbb{M}_k(t)$.*

Proof. We get the equivalence between the first two assertions by expanding:

$$\|A_i(t) - A\|^2 = \frac{1}{2} [\text{tr}(A_i(t)^T A_i(t)) + \text{tr}(A^T A)] - 2A \cdot A_i(t) = 3 - 2A \cdot A_i(t),$$

since A and $A_i(t)$ are both orthogonal matrices. So minimising the weighted sum of the squares distances amounts to maximising inner product of A and the weighted sum \mathbb{M}_k of the matrices A_i given by (3.8).

Therefore if $\det \mathbb{M}_k > 0$, and A is the polar decomposition of \mathbb{M}_k , we immediately get that $\det A > 0$, hence $A \in \text{SO}(3)$. We know that S can be diagonalised in an orthogonal basis: $S = P^T D P$ with $P^T P = \text{Id}$ and D is a diagonal matrix with positive diagonal elements $\lambda_1, \lambda_2, \lambda_3$. Now if $B \in \text{SO}(3)$ maximises $\frac{1}{2} \text{tr}(B^T \mathbb{M}_k)$ among all matrices in $\text{SO}(3)$, then it maximises $\text{tr}(B^T A P^T D P) = \text{tr}(P B^T A P^T D)$. So the matrix $\bar{B} = P B^T A P^T$ maximises $\text{tr}(\bar{B} D) = \lambda_1 \bar{b}_{11} + \lambda_2 \bar{b}_{22} + \lambda_3 \bar{b}_{33}$ among the elements of $\text{SO}(3)$ (the map $B \mapsto P B^T A P^T$ is a one-to-one correspondence between $\text{SO}(3)$ and itself). But since \bar{B} is an orthogonal matrix, all its column vectors are unit vectors, and so $\bar{b}_{ii} \leq 1$, with equality for $i = 1, 2$ and 3 if and only if $\bar{B} = \text{Id}$, that is to say $P B^T A P^T = \text{Id}$, which is exactly $B = A$. \square

We denote by $\text{PD}(\mathbb{M}_k) \in O(3)$ the corresponding orthogonal matrix coming from the polar decomposition of \mathbb{M}_k .

We now have two choices for the evolution of A_k . We can use the second point of Proposition 3.1 and follow the gradient of the function to maximise:

$$\frac{dA_k(t)}{dt} = \nu \nabla_A (\mathbb{M}_k \cdot A) |_{A=A_k} = \nu P_{T_{A_k}} \mathbb{M}_k \quad (3.9)$$

(see (A.2) for the last computation, $P_{T_{A_k}}$ is the projection on the tangent space, this way the solution of the equation stays in $\text{SO}(3)$).

Or we can directly relax to the polar decomposition $\text{PD}(\mathbb{M}_k)$, in the same manner:

$$\frac{dA_k(t)}{dt} = \nu P_{T_{A_k}} (\text{PD}(\mathbb{M}_k)).$$

We can actually see that the trajectory of this last equation, when $\text{PD}(\mathbb{M}_k)$ belongs to $\text{SO}(3)$ and does not depend on t , is exactly following a geodesic (see Proposition A.4). Therefore in this paper we will focus on this type of coordination. The positive coefficient ν gives the intensity of coordination, in the following we will assume that it is a function of the distance between A_k and $\text{PD}(\mathbb{M}_k)$ (the angle of the rotation $A_k^T \text{PD}(\mathbb{M}_k)$), which is equivalent to say that ν depends on $A_k \cdot \text{PD}(\mathbb{M}_k)$.

Remark 3.1. Some comments:

- (1) One could have used the Gram–Schmidt orthogonalisation instead of the polar decomposition, but it depends on the order in which the vector basis is taken (for instance if we start with \mathbf{e}_1 , it would define the first vector as the average of all the directions of displacement, independently of how the other vectors of the body attitudes of the individuals are distributed). The polar decomposition gives a more canonical way to do this.
- (2) We expect that the orthogonal matrix coming from the polar decomposition of \mathbb{M}_k belongs in fact to $\text{SO}(3)$. First, we notice that $O(3)$ is formed by two

disconnected components: $\text{SO}(3)$ and the other component formed by the matrices with determinant -1 . We assume that the motion of the agents is smooth enough so that the average \mathbb{M}_k stays “close” to $\text{SO}(3)$ and that, in particular, $\det(\mathbb{M}_k) > 0$.

A simple example is when we only average two different matrices A_1 and A_2 of $\text{SO}(3)$. We then have $\mathbb{M} = \frac{1}{2}(A_1 + A_2)$. If we write $A_1 A_2^T = \exp(\theta[\mathbf{n}]_\times)$ thanks to Rodrigues’ formula (3.4) and we define $A = A_2 \exp(\frac{1}{2}\theta[\mathbf{n}]_\times)$, we get that $A_1 = A \exp(\frac{1}{2}\theta[\mathbf{n}]_\times)$ and so

$$\begin{aligned}\mathbb{M} &= A \frac{1}{2} \left(\exp\left(\frac{1}{2}\theta[\mathbf{n}]_\times\right) + \exp\left(-\frac{1}{2}\theta[\mathbf{n}]_\times\right) \right) \\ &= A \left(\cos \frac{\theta}{2} \text{Id} + \left(1 - \cos \frac{\theta}{2}\right) \mathbf{n} \otimes \mathbf{n} \right),\end{aligned}$$

thanks to Rodrigues’ formula (3.3) and to (3.5). Since the matrix $S = \cos \frac{\theta}{2} \text{Id} + (1 - \cos \frac{\theta}{2}) \mathbf{n} \otimes \mathbf{n}$ is a positive-definite symmetric matrix as soon as $\theta \in [0, \pi)$, we have that $\det(\mathbb{M}) > 0$. The polar decomposition of \mathbb{M} is then A , which is the midpoint of the geodesic joining A_1 to A_2 (which corresponds to the curve $t \in [0, \theta] \mapsto A_1 \exp(t[\mathbf{n}]_\times)$).

As soon as we average more than two matrices, there exist cases for which $\det(\mathbb{M}) < 0$: for instance if we take

$$A_1 = \begin{pmatrix} 1 & 0 & 0 \\ 0 & -1 & 0 \\ 0 & 0 & -1 \end{pmatrix}, \quad A_2 = \begin{pmatrix} -1 & 0 & 0 \\ 0 & 1 & 0 \\ 0 & 0 & -1 \end{pmatrix}, \quad A_3 = \begin{pmatrix} -1 & 0 & 0 \\ 0 & -1 & 0 \\ 0 & 0 & 1 \end{pmatrix},$$

we have $\mathbb{M} = \frac{1}{3}(A_1 + A_2 + A_3) = -\frac{1}{3}\text{Id}$.

Noise term. Agents make errors when trying to coordinate their body attitude with that of their neighbours. This is represented in the equation of A_k by a noise term: $2\sqrt{D}dW_t^k$ where $D > 0$ and $W_t^k = (W_t^{k,i,j})_{i,j=1,2,3}$ are independent Gaussian distributions (Brownian motion).

From all these considerations, we obtain the IBM:

$$d\mathbf{X}_k(t) = v_0 A_k(t) \mathbf{e}_1 dt, \tag{3.10}$$

$$dA_k(t) = P_{T_{A_k}} \circ [\nu(\text{PD}(\mathbb{M}_k) \cdot A_k) \text{PD}(\mathbb{M}_k) dt + 2\sqrt{D}dW_t^k], \tag{3.11}$$

where the SDE is in Stratonovich sense (see Ref. 32). The projection $P_{T_{A_k}}$ and the fact that we consider the SDE in Stratonovich sense ensure that the solution $A_k(t)$ stays in $\text{SO}(3)$. The normalisation constant $2\sqrt{D}$ ensures that the diffusion coefficient is exactly D : the law p of the underlying process given by $dA_k = 2\sqrt{D}P_{T_{A_k}} \circ dW_t^k$ satisfies $\partial_t p = D\Delta_A p$ where $\Delta_A = \nabla_A \cdot \nabla_A$ is the Laplace–Beltrami operator on $\text{SO}(3)$. Notice the factor $2\sqrt{D}$ instead of the usual $\sqrt{2D}$ which

is encountered when considering diffusion process on manifolds isometrically embedded in the Euclidean space \mathbb{R}^n , because we are here considering $\text{SO}(3)$ embedded in \mathcal{M} (isomorphic to \mathbb{R}^9), but with the metric (3.6), which corresponds to the canonical metric of \mathbb{R}^9 divided by a factor 2. We refer to Ref. 37 for more insight on such stochastic processes on manifolds.

3.2. Mean-field limit

We assume that the kernel of influence K is Lipschitz, bounded, with the following properties:

$$K = K(|x|) \geq 0, \quad \int_{\mathbb{R}^3} K(|x|) dx = 1, \quad \int_{\mathbb{R}^3} |x|^2 K(|x|) dx < \infty. \quad (3.12)$$

In Ref. 7 the mean-field limit is proven for the Vicsek model. Using the techniques there it is straightforward to see that for

$$\mathbb{M}(x, t) := \frac{1}{N} \sum_{i=1}^N K(\mathbf{X}_i - x) A_i,$$

the law $f^N = f^N(x, A, t)$ of the empirical measure associated to the Stratonovich SDE:

$$d\mathbf{X}_{\mathbf{k}}(t) = v_0 A_{\mathbf{k}}(t) \mathbf{e}_1 dt, \quad (3.13)$$

$$dA_{\mathbf{k}}(t) = P_{T_{A_{\mathbf{k}}}} \circ [\nu(\mathbb{M}(\mathbf{X}_{\mathbf{k}}, t) \cdot A_{\mathbf{k}}) \mathbb{M}(\mathbf{X}_{\mathbf{k}}, t) dt + 2\sqrt{D} dW_t^k], \quad (3.14)$$

converges weakly $f^N \rightarrow f$ as $N \rightarrow \infty$. The limit satisfies the kinetic equation:

$$\partial_t f + v_0 A \mathbf{e}_1 \cdot \nabla_x f = D \Delta_A f - \nabla_A \cdot (F[f] f),$$

with

$$F[f] := \nu(\mathbb{M}_f \cdot A) P_{T_A}(\mathbb{M}_f),$$

$$\mathbb{M}_f = \int_{\mathbb{R}^3 \times \text{SO}(3)} K(x - x') f(x', A', t) A' dA' dx'.$$

The equations we are dealing with (3.10)–(3.11), since we consider the polar decomposition of the averaged body attitude \mathbb{M}_k , are slightly different from (3.13)–(3.14), which would correspond to the modelling point of view of Eq. (3.9). As a consequence, the corresponding coefficient of the SDE is not Lipschitz anymore and the known results for existence of solutions and mean-field limit (see Theorem 1.4 in Ref. 43) fail. More precisely, the problem arises when dealing with matrices with determinant zero; the orthogonal matrix of the polar decomposition is not uniquely defined for matrices with determinant zero and, otherwise, $\text{PD}(\mathbb{M}_k) = \mathbb{M}_k (\sqrt{\mathbb{M}_k^T \mathbb{M}_k})^{-1}$ (Lemma 3.1).

A complete proof of the previous results in the case of Eqs. (3.10)–(3.11) would involve proving that solutions to the equations stay away from the singular case $\det(\mathbb{M}_k) = 0$. This is an assumption that we make on the individual-based

model (see the second point of Remark 3.1). This kind of analysis has been done for the Vicsek model (explained in Sec. 1) in Ref. 28 where the authors prove global well-posedness for the kinetic equation in the spatially homogeneous case.

In our case one expects the following to hold.

Proposition 3.2. (Formal) *When the number of agents in (3.10)–(3.11) $N \rightarrow \infty$, its corresponding empirical distribution*

$$f^N(x, A, t) = \frac{1}{N} \sum_{k=1}^N \delta_{(X_k(t), A_k(t))},$$

converges weakly to $f = f(x, A, t)$, $(x, A, t) \in \mathbb{R}^3 \times \text{SO}(3) \times [0, \infty)$ satisfying

$$\begin{aligned} \partial_t f + v_0 A \mathbf{e}_1 \cdot \nabla_x f &= D \Delta_A f - \nabla_A \cdot (f F[f]), \\ F[f] &:= \nu P_{T_A}(\bar{\mathbb{M}}[f]), \\ \bar{\mathbb{M}}[f] &= \text{PD}(\mathbb{M}[f]), \\ \mathbb{M}[f](x, t) &:= \int_{\mathbb{R}^3 \times \text{SO}(3)} K(x - x') f(x', A', t) A' dA' dx', \end{aligned} \tag{3.15}$$

where $\text{PD}(\mathbb{M}[f])$ corresponds to the orthogonal matrix obtained on the polar decomposition of $\mathbb{M}[f]$ (see Lemma 3.1); and $\nu = \nu(\bar{\mathbb{M}}[f] \cdot A)$.

4. Hydrodynamic Limit

The goal of this section will be to derive the macroscopic equations (Theorem 4.1). From now on, we consider the kinetic equation given in (3.15).

4.1. Scaling and expansion

We express the kinetic equation (3.15) in dimensionless variables. Let ν_0 be the typical interaction frequency scale so that $\nu(\bar{A} \cdot A) = \nu_0 \nu'(\bar{A} \cdot A)$ with $\nu'(\bar{A} \cdot A) = \mathcal{O}(1)$. We introduce also the typical time and space scales t_0, x_0 such that $t_0 = \nu_0^{-1}$ and $x_0 = v_0 t_0$; the associated variables will be $t' = t/t_0$ and $x' = x/x_0$. Consider the dimensionless diffusion coefficient $d = D/\nu_0$ and the rescaled influence kernel $K'(|x'|) = K(x_0|x'|)$. Skipping the primes we get:

$$\begin{aligned} \partial_t f + A \mathbf{e}_1 \cdot \nabla_x f &= d \Delta_A f - \nabla_A \cdot (f F[f]), \\ F[f] &:= \nu(\bar{\mathbb{M}}[f] \cdot A) P_{T_A}(\bar{\mathbb{M}}[f]), \\ \bar{\mathbb{M}}[f] &= \text{PD}(\mathbb{M}[f]), \quad \mathbb{M}[f](x, t) := \int_{\mathbb{R}^3 \times \text{SO}(3)} K(x - x') f(x', A', t) A' dA' dx'. \end{aligned}$$

Here d, ν and K are assumed to be of order 1.

Remark 4.1. Notice in particular that before and after scaling the ratio

$$\frac{\nu}{D} = \frac{\nu'}{d}$$

remains the same.

Now, to carry out the macroscopic limit we rescale the space and time variables by setting $\tilde{t} = \varepsilon t$, $\tilde{x} = \varepsilon x$ to obtain (skipping the tildes):

$$\partial_t f^\varepsilon + A \mathbf{e}_1 \cdot \nabla_x f^\varepsilon = \frac{1}{\varepsilon} (d \Delta_A f^\varepsilon - \nabla_A \cdot (f^\varepsilon F^\varepsilon[f^\varepsilon])),$$

$$F^\varepsilon[f] := \nu(\bar{\mathbb{M}}^\varepsilon[f] \cdot A) P_{T_A}(\bar{\mathbb{M}}^\varepsilon[f^\varepsilon]),$$

$$\bar{\mathbb{M}}^\varepsilon[f] = \text{PD}(\mathbb{M}^\varepsilon[f]), \quad \mathbb{M}^\varepsilon[f](x, t) := \int_{\mathbb{R}^3 \times \text{SO}(3)} K\left(\frac{x - x'}{\varepsilon}\right) f(x', A', t) A' dA' dx'.$$

Lemma 4.1. *Assuming that f is sufficiently smooth (with bounded derivatives), we have the expansion*

$$\bar{\mathbb{M}}^\varepsilon[f](x, t) = \Lambda[f](x, t) + \mathcal{O}(\varepsilon^2),$$

where

$$\Lambda[f](x, t) = \text{PD}(\lambda[f]) \quad \text{and} \quad \lambda[f] = \int_{\text{SO}(3)} A' f(x, A', t) dA'.$$

Proof. This is obtained by performing the change of variable $x' = x + \varepsilon \xi$ in the definition of $\mathbb{M}^\varepsilon[f]$ and using a Taylor expansion of $f(x + \varepsilon \xi, A', t)$ with respect to ε . We use that K is isotropic and with bounded second moment by assumption (see Eq. (3.12)). \square

From the lemma, we rewrite:

$$\begin{aligned} \partial_t f^\varepsilon + A \mathbf{e}_1 \cdot \nabla_x f^\varepsilon &= \frac{1}{\varepsilon} Q(f^\varepsilon) + \mathcal{O}(\varepsilon), \\ F_0[f] &:= \nu(\Lambda[f] \cdot A) P_{T_A}(\Lambda[f]), \\ \Lambda[f] &= \text{PD}(\lambda[f]), \quad \lambda[f](x, t) := \int_{\text{SO}(3)} f(x, A', t) A' dA', \\ Q(f) &:= d \Delta_A f - \nabla_A \cdot (f F_0[f]). \end{aligned} \tag{4.1}$$

$\Lambda[f]$, $Q(f)$ and $F_0[f]$ are nonlinear operators of f , which only acts on the attitude variable A .

4.2. Preliminaries: Differential calculus in $\text{SO}(3)$

In the sequel we will use the volume form, the gradient and divergence in $\text{SO}(3)$ expressed in the Euler axis-angle coordinates (θ, \mathbf{n}) (explained at the beginning of Sec. 3). In this section, we give their explicit forms; the proofs are in Appendix A.

Proposition 4.1. (The gradient in $\text{SO}(3)$) *Let $f : \text{SO}(3) \rightarrow \mathbb{R}$ be a smooth scalar function. If $\bar{f}(\theta, \mathbf{n}) = f(A(\theta, \mathbf{n}))$ is the expression of f in the Euler axis-angle coordinates by Rodrigues' formula (3.3), we have then*

$$\nabla_A f = \partial_\theta \bar{f} A[\mathbf{n}]_\times + \frac{1}{2 \sin(\theta/2)} A(\cos(\theta/2) [\nabla_{\mathbf{n}} \bar{f}]_\times + \sin(\theta/2) [\mathbf{n} \times \nabla_{\mathbf{n}} \bar{f}]_\times), \quad (4.2)$$

where $A = A(\theta, \mathbf{n})$ and $\nabla_{\mathbf{n}}$ is the gradient on the sphere S^2 .

The volume form in $\text{SO}(3)$ is left invariant (it is the Haar measure), due to the fact that the inner product in \mathcal{M} is also left invariant: $A \cdot B = \frac{1}{2} \text{tr}(A^T B) = \Lambda A \cdot \Lambda B$ when $\Lambda \in \text{SO}(3)$. We give its expression in the Euler axis-angle coordinates (θ, \mathbf{n}) .

Lemma 4.2. (Decomposition of the volume form in $\text{SO}(3)$) *If $\bar{f}(\theta, \mathbf{n}) = f(A(\theta, \mathbf{n}))$ is the expression of f in the Euler axis-angle coordinates by Rodrigues' formula (3.3), we have*

$$\int_{\text{SO}(3)} f(A) dA = \int_0^\pi W(\theta) \int_{S^2} \bar{f}(\theta, \mathbf{n}) d\mathbf{n} d\theta,$$

where $d\mathbf{n}$ is the Lebesgue measure on the sphere S^2 , normalised to be a probability measure, and

$$W(\theta) = \frac{2}{\pi} \sin^2(\theta/2). \quad (4.3)$$

We have seen in Proposition 4.1 that the gradient is decomposed in the basis

$$\{A[\mathbf{n}]_\times, A[\nabla_{\mathbf{n}} \bar{f}]_\times, A[\mathbf{n} \times \nabla_{\mathbf{n}} \bar{f}]_\times\},$$

which are three orthogonal vectors of T_A (by Proposition A.2).

More generally if $B \in T_A$ for $A = A(\theta, \mathbf{n}) \in \text{SO}(3)$, then B is of the form AH with H antisymmetric, so $H = [\mathbf{u}]_\times$ for some $\mathbf{u} \in \mathbb{R}^3$. Decomposing \mathbf{u} on \mathbf{n} and its orthogonal, we get that there exist $\mathbf{v} \perp \mathbf{n}$ and $b \in \mathbb{R}$ such that $B = bA[\mathbf{n}]_\times + A[\mathbf{v}(\theta, \mathbf{n})]_\times$. Expressing B in this form, we compute the divergence in $\text{SO}(3)$.

Proposition 4.2. (The divergence in $\text{SO}(3)$) *Consider $B : \text{SO}_3 \rightarrow T(\text{SO}(3))$ a smooth function (so that $B(A) \in T_A$ for all $A \in \text{SO}(3)$), and suppose that*

$$B(A(\theta, \mathbf{n})) = b(\theta, \mathbf{n}) A[\mathbf{n}]_\times + A[\mathbf{v}(\theta, \mathbf{n})]_\times,$$

for some smooth function b and smooth vector function \mathbf{v} such that $\mathbf{v}(\theta, \mathbf{n}) \perp \mathbf{n}$. Then

$$\begin{aligned} \nabla_A \cdot B &= \frac{1}{\sin^2(\theta/2)} \partial_\theta (\sin^2(\theta/2) b(\theta, \mathbf{n})) \\ &\quad + \frac{1}{2 \sin(\theta/2)} \nabla_{\mathbf{n}} \cdot (\mathbf{v}(\theta, \mathbf{n}) \cos(\theta/2) + (\mathbf{v}(\theta, \mathbf{n}) \times \mathbf{n}) \sin(\theta/2)). \end{aligned}$$

Now we can compute the Laplacian in $\text{SO}(3)$.

Corollary 4.1. *The Laplacian in $\text{SO}(3)$ can be expressed as*

$$\Delta_A f = \frac{1}{\sin^2(\theta/2)} \partial_\theta (\sin^2(\theta/2) \partial_\theta \tilde{f}) + \frac{1}{4 \sin^2(\theta/2)} \Delta_{\mathbf{n}} \tilde{f},$$

where $\Delta_{\mathbf{n}}$ is the Laplacian on the sphere S^2 and $f(A) = f(A(\theta, \mathbf{n})) = \tilde{f}(\theta, \mathbf{n})$.

Proof. Let $B(\theta, \mathbf{n}) := \nabla_A f(A(\theta, \mathbf{n})) \in T_A$. Then, using the notations of Proposition 4.2 and the result of Proposition 4.1, we have that:

$$b = \partial_\theta \tilde{f},$$

$$\mathbf{v} = \frac{1}{2 \sin(\theta/2)} (\cos(\theta/2) \nabla_{\mathbf{n}} \tilde{f} + \sin(\theta/2) (\mathbf{n} \times \nabla_{\mathbf{n}} \tilde{f})),$$

from here we just need to apply Proposition 4.2 knowing that $(\mathbf{n} \times \nabla_{\mathbf{n}} \tilde{f}) \times \mathbf{n} = \nabla_{\mathbf{n}} \tilde{f}$ since $\nabla_{\mathbf{n}} \tilde{f}$ is orthogonal to \mathbf{n} . \square

4.3. Equilibrium solutions and Fokker–Planck formulation

We define a generalisation of the von Mises distributions on $\text{SO}(3)$ by

$$M_\Lambda(A) = \frac{1}{Z} \exp\left(\frac{\sigma(A \cdot \Lambda)}{d}\right), \quad \int_{\text{SO}(3)} M_\Lambda(A) dA = 1, \quad \Lambda \in \text{SO}(3), \quad (4.4)$$

where $Z = Z(\nu, d)$ is a normalising constant and $\sigma = \sigma(\mu)$ is such that $(d/d\mu)\sigma = \nu(\mu)$. Observe that $Z < \infty$ is independent of Λ since the volume form on $\text{SO}(3)$ is left invariant. Therefore we have

$$\begin{aligned} Z &= \int_{\text{SO}(3)} \exp(d^{-1} \sigma(A \cdot \Lambda)) dA = \int_{\text{SO}(3)} \exp(d^{-1} \sigma(\Lambda^T A \cdot \text{Id})) dA \\ &= \int_{\text{SO}(3)} \exp(d^{-1} \sigma(A \cdot \text{Id})) dA, \end{aligned}$$

and we also obtain that $M_\Lambda(A)$ is actually $M_{\text{Id}}(\Lambda^T A)$.

We are now ready to describe the properties of Q in terms of these generalised von Mises distributions.

Lemma 4.3. (Properties of Q) *The following holds:*

(i) *The operator Q can be written as*

$$Q(f) = d \nabla_A \cdot \left[M_{\Lambda[f]} \nabla_A \left(\frac{f}{M_{\Lambda[f]}} \right) \right]$$

and we have

$$H(f) := \int_{\text{SO}(3)} Q(f) \frac{f}{M_{\Lambda[f]}} dA = -d \int_{\text{SO}(3)} M_{\Lambda[f]} \left| \nabla_A \left(\frac{f}{M_{\Lambda[f]}} \right) \right|^2 dA \leq 0. \quad (4.5)$$

- (ii) The equilibria, i.e. the functions $f = f(x, A, t)$ such that $Q(f) = 0$ form a four-dimensional manifold \mathcal{E} given by

$$\mathcal{E} = \{\rho M_\Lambda(A) \mid \rho > 0, \Lambda \in \text{SO}(3)\},$$

where ρ is the total mass while Λ is mean body attitude of $\rho M_\Lambda(A)$, i.e.:

$$\rho = \int_{\text{SO}(3)} \rho M_\Lambda(A) dA,$$

$$\Lambda = \Lambda[\rho M_\Lambda].$$

Furthermore, $H(f) = 0$ if and only if $f = \rho M_\Lambda$ for arbitrary $\rho \in \mathbb{R}_+$ and $\Lambda \in \text{SO}(3)$.

To prove Lemma 4.3 we require the following one, which is of independent interest and for which we introduce the following notation: for any scalar function $g : (0, \pi) \rightarrow \mathbb{R}$ and a given integrable scalar function $h : (0, \pi) \rightarrow \mathbb{R}$ which remains positive (or negative) on $(0, \pi)$, we define

$$\langle g(\theta) \rangle_{h(\theta)} := \int_0^\pi g(\theta) \frac{h(\theta)}{\int_0^\pi h(\theta') d\theta'} d\theta. \quad (4.6)$$

Lemma 4.4. (Consistency relation for the “flux”)

$$\lambda[M_{\Lambda_0}] = c_1 \Lambda_0,$$

where $c_1 \in (0, 1)$ is equal to

$$c_1 = \frac{2}{3} \left\langle \frac{1}{2} + \cos \theta \right\rangle_{m(\theta) \sin^2(\theta/2)} \quad (4.7)$$

for

$$m(\theta) = \exp \left(d^{-1} \sigma \left(\frac{1}{2} + \cos \theta \right) \right). \quad (4.8)$$

Proof. Using the fact that the measure on $\text{SO}(3)$ is left invariant, we obtain

$$\begin{aligned} \lambda[M_{\Lambda_0}] &= \frac{1}{Z} \int_{\text{SO}(3)} A \exp(d^{-1} \sigma((A \cdot \Lambda_0))) dA \\ &= \frac{\Lambda_0}{Z} \int_{\text{SO}(3)} \Lambda_0^T A \exp \left(d^{-1} \sigma \left(\frac{1}{2} \text{tr}(\Lambda_0^T A) \right) \right) dA \\ &= \frac{\Lambda_0}{Z} \int_{\text{SO}(3)} B \exp \left(d^{-1} \sigma \left(\frac{1}{2} \text{tr}(B) \right) \right) dB. \end{aligned}$$

We now write $B = \text{Id} + \sin \theta [\mathbf{n}]_{\times} + (1 - \cos \theta) [\mathbf{n}]_{\times}^2$ thanks to Rodrigues' formula (3.3). Therefore, using Lemma 4.2, we get

$$\begin{aligned} \lambda[M_{\Lambda_0}] &= \Lambda_0 \frac{\int_{\text{SO}(3)} B \exp(d^{-1} \sigma(\frac{1}{2} \text{tr}(B))) dB}{\int_{\text{SO}(3)} \exp(d^{-1} \sigma(\frac{1}{2} \text{tr}(B))) dB} \\ &= \Lambda_0 \frac{\int_0^\pi \sin^2(\theta/2) \exp(d^{-1} \sigma(\frac{1}{2} + \cos \theta)) \times (\int_{S^2} (\text{Id} + \sin \theta [\mathbf{n}]_{\times} + (1 - \cos \theta) [\mathbf{n}]_{\times}^2) d\mathbf{n}) d\theta}{\int_0^\pi \sin^2(\theta/2) \exp(d^{-1} \sigma(\frac{1}{2} + \cos \theta)) d\theta}. \end{aligned}$$

Next, we see that since the function $\mathbf{n} \mapsto [\mathbf{n}]_{\times}$ is odd, we have $\int_{S^2} [\mathbf{n}]_{\times} d\mathbf{n} = 0$. We also have (see (3.5)) that $[\mathbf{n}]_{\times}^2 = \mathbf{n} \otimes \mathbf{n} - \text{Id}$. Since we know that $\int_{S^2} \mathbf{n} \otimes \mathbf{n} d\mathbf{n} = \frac{1}{3} \text{Id}$ (by invariance by rotation), it is easy to see that the integral in S^2 has to be proportional to Id , the coefficient is given by computing the trace), we get that

$$\begin{aligned} \lambda[M_{\Lambda_0}] &= \Lambda_0 \frac{\int_{\text{SO}(3)} B \exp(d^{-1} \sigma(\frac{1}{2} \text{tr}(B))) dB}{\int_{\text{SO}(3)} \exp(d^{-1} \sigma(\frac{1}{2} \text{tr}(B))) dB} \\ &= \Lambda_0 \frac{\int_0^\pi \sin^2(\theta/2) \exp(d^{-1} \sigma(\frac{1}{2} + \cos \theta)) (\text{Id} + (1 - \cos \theta) (\frac{1}{3} - 1) \text{Id}) d\theta}{\int_0^\pi \sin^2(\theta/2) \exp(d^{-1} \sigma(\frac{1}{2} + \cos \theta)) d\theta} \\ &= \frac{\int_0^\pi \frac{2}{3} (\frac{1}{2} + \cos \theta) \sin^2(\theta/2) \exp(d^{-1} \sigma(\frac{1}{2} + \cos \theta)) d\theta}{\int_0^\pi \sin^2(\theta/2) \exp(d^{-1} \sigma(\frac{1}{2} + \cos \theta)) d\theta} \Lambda_0 = c_1 \Lambda_0, \end{aligned}$$

which gives the formula (4.7) for c_1 .

It remains to prove that $c_1 \in (0, 1)$. We have that c_1 is the average of $\frac{2}{3}(\frac{1}{2} + \cos \theta)$ for the probability measure on $(0, \pi)$ proportional to $\sin^2(\theta/2) \exp(d^{-1} \sigma(\frac{1}{2} + \cos \theta))$. Since we have $\frac{2}{3}(\frac{1}{2} + \cos \theta) \leq 1$ with equality only for $\theta = 0$, we immediately get that $c_1 < 1$. To prove the positivity, we remark that the function in the exponent $\theta \mapsto d^{-1} \sigma(\frac{1}{2} + \cos \theta)$ is strictly decreasing for $\theta \in (0, \pi)$ (since $\nu > 0$ is the derivative of σ), so we obtain that $\sigma(\frac{1}{2} + \cos \theta) > \sigma(\frac{1}{2} + \cos \frac{2\pi}{3}) = \sigma(0)$ for $\theta \in (0, \frac{2\pi}{3})$. Therefore, for $\theta \in (0, \frac{2\pi}{3})$,

$$\left(\frac{1}{2} + \cos \theta\right) \exp\left(d^{-1} \sigma\left(\frac{1}{2} + \cos \theta\right)\right) > \left(\frac{1}{2} + \cos \theta\right) \exp(d^{-1} \sigma(0)),$$

since $\frac{1}{2} + \cos \theta > 0$. When $\theta \in (\frac{2\pi}{3}, \pi)$, we have exactly the same inequality above since we have $\frac{1}{2} + \cos \theta < 0$. Therefore we get

$$c_1 > \frac{\int_0^\pi \frac{2}{3} (\frac{1}{2} + \cos \theta) \sin^2(\theta/2) \exp(d^{-1} \sigma(0)) d\theta}{\int_0^\pi \sin^2(\theta/2) \exp(d^{-1} \sigma(\frac{1}{2} + \cos \theta)) d\theta} = 0,$$

since

$$\begin{aligned} \int_0^\pi \left(\frac{1}{2} + \cos \theta \right) \sin^2(\theta/2) d\theta &= \int_0^\pi \left(\frac{1}{2} + \cos \theta \right) \left(\frac{1}{2} - \frac{1}{2} \cos \theta \right) d\theta \\ &= \frac{\pi}{4} - \frac{1}{2} \int_0^\pi \cos^2 \theta d\theta = 0. \end{aligned} \quad \square$$

Proof of Lemma 4.3. We follow the structure of the analogous proof in Ref. 24:

- (i) To prove the first identity we have that (see expression (A.2)):

$$\begin{aligned} \nabla_A (\ln M_{\Lambda[f]}) &= d^{-1} \nabla_A (\sigma(A \cdot \Lambda[f])) \\ &= d^{-1} \nu(A \cdot \Lambda[f]) P_{T_A}(\Lambda[f]) \\ &= d^{-1} F_0[f] \end{aligned}$$

and so

$$\begin{aligned} d\nabla_A \cdot \left[M_{\Lambda[f]} \nabla_A \left(\frac{f}{M_{\Lambda[f]}} \right) \right] &= d\nabla_A \cdot [\nabla_A f - f \nabla_A (\ln(M_{\Lambda[f]}))] \\ &= d\nabla_A f - \nabla_A \cdot (f F_0[f]). \end{aligned}$$

Inequality (4.5) follows from this last expression and the Stokes theorem in $\text{SO}(3)$.

- (ii) From the inequality (4.5) we have that if $Q(f) = 0$, then $\frac{f}{M_{\Lambda[f]}}$ is a constant that we denote by ρ (which is positive since f and $M_{\Lambda[f]}$ are positive). Conversely, if $f = \rho M_\Lambda$ then

$$\lambda[\rho M_\Lambda] = \int_{\text{SO}(3)} \rho M_\Lambda(A) A dA = \rho c_1 \Lambda$$

by Lemma 4.4. Now, by uniqueness of the polar decomposition and since $\rho c_1 \text{Id}$ is a symmetric positive-definite matrix, we have that $\Lambda[\rho M_\Lambda] = \Lambda$. \square

Let us describe the behaviour of these equilibrium distributions for small and large noise intensities. We have that for any function g , the average $\langle g(\frac{1}{2} + \cos \theta) \rangle_{m(\theta) \sin^2(\theta/2)}$ is the average of $g(A \cdot \Lambda)$ with respect to the probability measure M_Λ (by left invariance, this is independent of Λ).

One can actually check that the probability measure M_Λ on $\text{SO}(3)$ converges in distribution to the uniform measure when $d \rightarrow \infty$ (by Taylor expansion) and it converges to a Dirac delta at matrix Λ when $d \rightarrow 0$ (this can be seen for M_{Id} thanks to the decomposition of the volume form and the Laplace method, since the maximum of $\sigma(\frac{1}{2} + \cos \theta)$ is reached only at $\theta = 0$ which corresponds to the identity matrix, and we then get the result for any Λ since $M_\Lambda(A) = M_{\text{Id}}(\Lambda^T A)$). So for small diffusion, at equilibrium, agents tend to adopt the same body attitude close to Λ .

With these asymptotic considerations, we have in particular the behaviour of c_1 :

$$c_1 \xrightarrow{d \rightarrow \infty} 0$$

and

$$c_1 \xrightarrow{d \rightarrow 0} 1.$$

4.4. Generalised collision invariants

To obtain the macroscopic equation, we start by looking for the conserved quantities of the kinetic equation: we want to find the functions $\psi = \psi(A)$ such that

$$\int_{\text{SO}(3)} Q(f) \psi dA = 0 \quad \text{for all } f.$$

By Lemma 4.3, this can be rewritten as

$$0 = - \int_{\text{SO}(3)} M_\Lambda[f] \nabla_A \left(\frac{f}{M_\Lambda[f]} \right) \cdot \nabla_A \psi dA.$$

This happens if $\nabla_A \psi \in T_A^\perp$ which holds true only if $\nabla_A \psi = 0$, implying that ψ is constant.

Consequently, our model has only one conserved quantity: the total mass. However the equilibria is four-dimensional (by Lemma 4.3). To obtain the macroscopic equations for Λ , *a priori* we would need three more conserved quantities. This problem is sorted out by using GCI a concept first introduced in Ref. 24.

4.4.1. Definition and existence of GCI

Define the operator

$$\mathcal{Q}(f, \Lambda_0) := \nabla_A \cdot \left(M_{\Lambda_0} \nabla_A \left(\frac{f}{M_{\Lambda_0}} \right) \right),$$

notice in particular that

$$Q(f) = \mathcal{Q}(f, \Lambda[f]).$$

Using this operator we define the following.

Definition 4.1. (GCI) For a given $\Lambda_0 \in \text{SO}(3)$ we say that a real-valued function $\psi : \text{SO}(3) \rightarrow \mathbb{R}$ is a GCI associated to Λ_0 , or for short $\psi \in \text{GCI}(\Lambda_0)$, if

$$\int_{\text{SO}(3)} \mathcal{Q}(f, \Lambda_0) \psi dA = 0 \quad \text{for all } f \text{ such that } P_{T_{\Lambda_0}}(\lambda[f]) = 0.$$

In particular, the result that we will use is

$$\psi \in \text{GCI}(\Lambda[f]) \Rightarrow \int_{\text{SO}(3)} Q(f) \psi dA = 0. \quad (4.9)$$

Indeed, since $\Lambda[f]$ is the polar decomposition of $\lambda[f]$, we have $\lambda[f] = \Lambda[f]S$, with S a symmetric matrix. Therefore (see Proposition A.2), we get that $\lambda[f]$ belongs to

the orthogonal of $T_{\Lambda[f]}$, so Definition 4.1 and the fact that $Q(f) = \mathcal{Q}(f, \Lambda[f])$ give us the property (4.9).

Definition 4.1 is equivalent to the following.

Proposition 4.3. *We have that $\psi \in \text{GCI}(\Lambda_0)$ if and only if:*

$$\text{there exists } B \in T_{\Lambda_0} \text{ such that } \nabla_A \cdot (M_{\Lambda_0} \nabla_A \psi) = B \cdot A M_{\Lambda_0}. \quad (4.10)$$

Proof. We denote by \mathcal{L} the linear operator $\mathcal{Q}(\cdot, \Lambda_0)$, and \mathcal{L}^* its adjoint. We have the following sequence of equivalences, starting from Definition 4.1:

$$\begin{aligned} \psi \in \text{GCI}(\Lambda_0) &\Leftrightarrow \int_{\text{SO}(3)} \psi \mathcal{L}(f) dA = 0 \quad \text{for all } f \text{ such that } P_{T_{\Lambda_0}}(\lambda[f]) = 0 \\ &\Leftrightarrow \int_{\text{SO}(3)} \mathcal{L}^*(\psi) f dA = 0 \quad \text{for all } f \text{ such that} \\ &\quad \int_{\text{SO}(3)} A f(A) dA \in (T_{\Lambda_0})^\perp \\ &\Leftrightarrow \int_{\text{SO}(3)} \mathcal{L}^*(\psi) f dA = 0 \quad \text{for all } f \text{ such that } \forall B \in T_{\Lambda_0}, \\ &\quad \int_{\text{SO}(3)} (B \cdot A) f(A) dA = 0 \\ &\Leftrightarrow \int_{\text{SO}(3)} \mathcal{L}^*(\psi) f dA = 0 \quad \text{for all } f \in F_{\Lambda_0}^\perp \\ &\Leftrightarrow \mathcal{L}^*(\psi) \in (F_{\Lambda_0}^\perp)^\perp, \end{aligned}$$

where

$$F_{\Lambda_0} := \{g : \text{SO}(3) \rightarrow \mathbb{R}, \text{ with } g(A) = (B \cdot A), \text{ for some } B \in T_{\Lambda_0}\},$$

and $F_{\Lambda_0}^\perp$ is the space orthogonal to F_{Λ_0} in L^2 . F_{Λ_0} is a vector space in L^2 -isomorphic to T_{Λ_0} and $(F_{\Lambda_0}^\perp)^\perp = F_{\Lambda_0}$ since F_{Λ_0} is closed (finite-dimensional). Therefore we get

$$\begin{aligned} \psi \in \text{GCI}(\Lambda_0) &\Leftrightarrow \mathcal{L}^*(\psi) \in F_{\Lambda_0} \\ &\Leftrightarrow \text{there exists } B \in T_{\Lambda_0} \text{ such that } \mathcal{L}^*(\psi)(A) = B \cdot A, \end{aligned}$$

which ends the proof since the expression of the adjoint is

$$\mathcal{L}^*(\psi) = \frac{1}{M_{\Lambda_0}} \nabla_A \cdot (M_{\Lambda_0} \nabla_A \psi). \quad \square$$

We prove the existence and uniqueness of the solution ψ satisfying Eq. (4.10) in the following.

Proposition 4.4. (Existence of the GCI) *For a given $B \in T_{\Lambda}$ fixed, there exists a unique (up to a constant) $\psi_B \in H^1(\text{SO}(3))$, satisfying the relation (4.10).*

Proof. We would like to apply the Lax–Milgram theorem to prove the existence of ψ in an appropriate functional space. For this, we rewrite the relation (4.10) weakly

$$a(\psi, \varphi) := \int_{\text{SO}(3)} M_{\Lambda_0} \nabla_A \psi \cdot \nabla_A \varphi dA = \int_{\text{SO}(3)} B \cdot P_{T_{\Lambda_0}}(A) M_{\Lambda_0} \varphi dA =: b(\varphi). \quad (4.11)$$

Our goal is to prove that there exists a unique $\psi \in H^1(\text{SO}(3))$ such that $a(\psi, \varphi) = b(\varphi)$ for all $\varphi \in H^1(\text{SO}(3))$.

To begin with we apply the Lax–Milgram theorem on the space

$$H_0^1(\text{SO}(3)) := \left\{ \varphi \in H^1 \left| \int_{\text{SO}(3)} \varphi dA = 0 \right. \right\}.$$

In this space the H^1 -norm and the H^1 -semi-norm are equivalent thanks to the Poincaré inequality, i.e. there exists $C > 0$ such that

$$\int_{\text{SO}(3)} |\nabla_A \varphi|^2 dA \geq C \int_{\text{SO}(3)} |\varphi|^2 dA \quad \text{for some } C > 0, \text{ for all } \varphi \in H_0^1(\text{SO}(3)).$$

Notice that the Poincaré inequality holds in $\text{SO}(3)$ because it is compact Riemannian manifold.¹² This gives us the coercivity estimate to apply the Lax–Milgram theorem. Hence, there exists a unique $\psi \in H_0^1(\text{SO}(3))$ such that $a(\psi, \varphi) = b(\varphi)$ for all $\varphi \in H_0^1(\text{SO}(3))$.

Now, define for a given $\varphi \in H^1(\text{SO}(3))$, $\varphi_0 := \varphi - \int_{\text{SO}(3)} \varphi dA \in H_0^1(\text{SO}(3))$. It holds that

$$a(\psi, \varphi) = a(\psi, \varphi_0) \quad \text{and} \quad b(\varphi) = b(\varphi_0),$$

since $b(1) = 0$ given that it has antisymmetric integrand. Hence, we obtain that there exists a unique $\psi \in H_0^1(\text{SO}(3))$ such that

$$a(\psi, \varphi) = b(\varphi) \quad \text{for all } \varphi \in H^1(\text{SO}(3)).$$

Suppose next, that there exists another solution $\bar{\psi} \in H^1(\text{SO}(3))$ to this problem, then the difference $\Psi = \psi - \bar{\psi}$ satisfies:

$$0 = a(\Psi, \varphi) = \int_{\text{SO}(3)} M_{\Lambda_0} \nabla_A \Psi \cdot \nabla_A \varphi dA \quad \text{for all } \varphi \in H^1(\text{SO}(3)).$$

Take in particular $\varphi = \Psi$, then

$$\int_{\text{SO}(3)} M_{\Lambda_0} |\nabla_A \Psi|^2 dA = 0.$$

Hence, $\Psi = c$ for some constant c , so all solutions are of the form $\psi + c$ where ψ is the unique solution satisfying $\int_{\text{SO}(3)} \psi dA = 0$. \square

By writing that

$$B \in T_{\Lambda_0} \quad \text{if and only if there exists } P \in \mathcal{A}, B = \Lambda_0 P, \quad (4.12)$$

with \mathcal{A} the set of antisymmetric matrices, we deduce the following.

Corollary 4.2. *For a given $\Lambda_0 \in \text{SO}(3)$, the set of GCIs associated to Λ_0 are*

$$\text{GCI}(\Lambda_0) = \text{span} \left\{ 1, \bigcup_{P \in \mathcal{A}} \psi_P^{\Lambda_0} \right\}$$

(where \mathcal{A} is the set of antisymmetric matrices) with $\psi_P^{\Lambda_0}$ the unique solution in $H_0^1(\text{SO}(3))$ of

$$a(\psi_P^{\Lambda_0}, \varphi) = b_P(\varphi) \quad \text{for all } \varphi \in H^1(\text{SO}(3)),$$

where a and b_P are defined by (4.11) with B substituted by $\Lambda_0 P$.

Notice that since the mapping $P \mapsto \psi_P^{\Lambda_0}$ is linear and injective from \mathcal{A} (of dimension 3) to $H_0^1(\text{SO}(3))$, the vector space $\text{GCI}(\Lambda_0)$ is of dimension 4.

4.4.2. The non-constant GCIs

From now on, we omit the subscript on Λ_0 , and we are interested in a simpler expression for ψ_P^Λ . Rewriting expression (4.10) using (4.12), for any given $P \in \mathcal{A}$ we want to find ψ such that

$$\nabla_A \cdot (M_\Lambda \nabla_A \psi) = (\Lambda P) \cdot A M_\Lambda = P \cdot (\Lambda^T A) M_\Lambda, \quad P \in \mathcal{A}. \quad (4.13)$$

Proposition 4.5. *Let $P \in \mathcal{A}$ and ψ be the solution of (4.13) belonging to $H_0^1(\text{SO}(3))$. If we denote $\bar{\psi}(B) := \psi(\Lambda B)$, then $\bar{\psi}$ is the unique solution in $H_0^1(\text{SO}(3))$ of*

$$\nabla_B \cdot (M_{\text{Id}}(B) \nabla_B \bar{\psi}) = P \cdot B M_{\text{Id}}(B). \quad (4.14)$$

Proof. Let $\psi(A) = \bar{\psi}(\Lambda^T A)$. Consider $A(\varepsilon)$ a differentiable curve in $\text{SO}(3)$ with

$$A(0) = A, \quad \left. \frac{dA(\varepsilon)}{d\varepsilon} \right|_{\varepsilon=0} = \delta_A \in T_A.$$

Then, by definition

$$\lim_{\varepsilon \rightarrow 0} \frac{\psi(A(\varepsilon)) - \psi(A)}{\varepsilon} = \nabla_A \psi(A) \cdot \delta_A,$$

and therefore we have that

$$\lim_{\varepsilon \rightarrow 0} \frac{\bar{\psi}(\Lambda^T A(\varepsilon)) - \bar{\psi}(\Lambda^T A)}{\varepsilon} = \nabla_B \bar{\psi}(\Lambda^T A) \cdot \Lambda^T \delta_A$$

since

$$\Lambda^T A(0) = \Lambda^T A, \quad \left. \frac{d}{d\varepsilon} \Lambda^T A(\varepsilon) \right|_{\varepsilon=0} = \Lambda^T \delta_A.$$

We conclude that

$$\nabla_A \psi(A) \cdot \delta_A = \nabla_B \bar{\psi}(\Lambda^T A) \cdot \Lambda^T \delta_A.$$

Now we check that

$$\begin{aligned}\frac{1}{2}\mathrm{tr}((\nabla_A\psi(A))^T\delta_A) &= \frac{1}{2}\mathrm{tr}((\nabla_B\bar{\psi}(\Lambda^TA))^T\Lambda^T\delta_A) \\ &= \frac{1}{2}\mathrm{tr}((\Lambda\nabla_B\bar{\psi}(\Lambda^TA))^T\delta_A),\end{aligned}$$

implying (since this is true for any $\delta_A \in T_A$) that

$$\nabla_A\psi(A) = \Lambda\nabla_B\bar{\psi}(\Lambda^TA).$$

Now to deal with the divergence term, we consider the variational formulation. Consider $\varphi \in H^1(\mathrm{SO}(3))$, then our equation is equivalent to

$$-\int_{\mathrm{SO}(3)} M_\Lambda(A)\nabla_A\psi(A) \cdot \nabla_A\varphi(A)dA = \int_{\mathrm{SO}(3)} P \cdot (\Lambda^TA)M_\Lambda(A)\varphi(A)dA,$$

for all $\varphi \in H^1(\mathrm{SO}(3))$. The left-hand side can be written as

$$\begin{aligned}&-\int_{\mathrm{SO}(3)} M_{\mathrm{Id}}(B)(\Lambda^TA) (\Lambda\nabla_B\bar{\psi}(\Lambda^TA)) \cdot (\Lambda\nabla_B\bar{\varphi}(\Lambda^TA))dA \\ &= -\int_{\mathrm{SO}(3)} M_{\mathrm{Id}}(B)\nabla_B\bar{\psi}(B) \cdot \nabla_B\bar{\varphi}(B)dB;\end{aligned}$$

and the right-hand side is equal to

$$\int_{\mathrm{SO}(3)} P \cdot BM_{\mathrm{Id}}(B)\bar{\varphi}(B)dB,$$

where we define analogously $\bar{\varphi}(B) = \varphi(\Lambda B)$. This concludes the proof. \square

Therefore it is enough to find the solution to (4.14). Inspired by Ref. 24 we make the ansatz:

$$\bar{\psi}(B) = P \cdot B\bar{\psi}_0 \left(\frac{1}{2}\mathrm{tr}(B) \right),$$

for some scalar function $\bar{\psi}_0$.

Proposition 4.6. (Non-constant GCI) *Let $P \in \mathcal{A}$, then the unique solution $\bar{\psi} \in H_0^1(\mathrm{SO}(3))$ of (4.14) is given by*

$$\bar{\psi}(B) = P \cdot B\bar{\psi}_0 \left(\frac{1}{2}\mathrm{tr}(B) \right), \quad (4.15)$$

where $\bar{\psi}_0$ is constructed as follows: let $\tilde{\psi}_0 : \mathbb{R} \rightarrow \mathbb{R}$ be the unique solution to

$$\frac{1}{\sin^2(\theta/2)}\partial_\theta(\sin^2(\theta/2)m(\theta)\partial_\theta(\sin\theta\tilde{\psi}_0)) - \frac{m(\theta)\sin\theta}{2\sin^2(\theta/2)}\tilde{\psi}_0 = \sin\theta m(\theta), \quad (4.16)$$

where $m(\theta) = M_{\mathrm{Id}}(B) = \exp(d^{-1}\sigma(\frac{1}{2} + \cos\theta))/Z$. Then

$$\tilde{\psi}_0(\theta) = \bar{\psi}_0 \left(\frac{1}{2}\mathrm{tr}(B) \right) \quad (4.17)$$

by the relation $\frac{1}{2}\text{tr}(B) = \frac{1}{2} + \cos\theta$. $\tilde{\psi}_0$ is 2π -periodic, even and negative (by the maximum principle).

Going back to the GCI $\psi(A)$, we can write it as

$$\psi(A) = P \cdot (\Lambda^T A) \bar{\psi}_0(\Lambda \cdot A). \quad (4.18)$$

Proof. Suppose that the solution is given by expression (4.15). We check that $\tilde{\psi}_0$ given by Eq. (4.17) satisfies Eq. (4.16) using the gradient and divergence in $\text{SO}(3)$ computed in Propositions 4.1 and 4.2. First notice that P is antisymmetric, thus if we write Rodrigues' formula (3.3) for $B(\theta, \mathbf{n})$, the symmetric part of $B(\theta, \mathbf{n})$ gives no contribution when computing $P \cdot B$ and we get

$$\bar{\psi}(B) = P \cdot B \bar{\psi}_0 \left(\frac{1}{2} \text{tr}(B) \right) = \sin \theta \tilde{\psi}_0(\theta) P \cdot [\mathbf{n}]_{\times} = \sin \theta \tilde{\psi}_0(\theta) (\mathbf{p} \cdot \mathbf{n}),$$

where the vector \mathbf{p} is such that $P = [\mathbf{p}]_{\times}$ and this leads to

$$\begin{aligned} \nabla_B \cdot (M_{\text{Id}}(B) \nabla_B \bar{\psi}) &= \frac{1}{\sin^2(\theta/2)} \partial_{\theta} (\sin^2(\theta/2) m(\theta) \partial_{\theta} (\sin \theta \tilde{\psi}_0(\theta))) (\mathbf{p} \cdot \mathbf{n}) \\ &\quad + \frac{m(\theta) \sin \theta}{4 \sin^2(\theta/2)} \tilde{\psi}_0(\theta) \Delta_n (\mathbf{p} \cdot \mathbf{n}). \end{aligned}$$

Using that the Laplacian in the sphere has the property

$$\Delta_n (\mathbf{p} \cdot \mathbf{n}) = -2(\mathbf{p} \cdot \mathbf{n})(\mathbf{p} \cdot \mathbf{n})$$

corresponds to the first spherical harmonic), we conclude that expression (4.16) is satisfied. In the computation we used the same procedure as for the proof of the expression of the Laplacian in $\text{SO}(3)$ (Corollary 4.1), but (using the same notations) we have taken $b(\theta, \mathbf{n}) = m(\theta) \partial_{\theta} (\sin \theta \tilde{\psi}_0(\theta)) (\mathbf{p} \cdot \mathbf{n})$.

To conclude the proof we just need to check that $\tilde{\psi}_0$ exists and corresponds to a function $\bar{\psi}$ in $H_0^1(\text{SO}(3))$. Using the expression of the volume form, since $\int_{S^2} \mathbf{p} \cdot \mathbf{n} d\mathbf{n} = 0$, we get that if ψ_0 is smooth, we have $\int_{\text{SO}(3)} \bar{\psi}(A) dA = 0$, and using the expression of the gradient, we get that

$$\begin{aligned} \int_{\text{SO}(3)} |\nabla \bar{\psi}(A)|^2 dA &= \frac{2}{\pi} \int_0^{\pi} \sin^2(\theta/2) |\partial_{\theta} (\sin \theta \tilde{\psi}_0(\theta))|^2 d\theta \int_{S^2} |\mathbf{p} \cdot \mathbf{n}|^2 d\mathbf{n} \\ &\quad + \frac{2}{\pi} \int_0^{\pi} \frac{1}{4} |\sin \theta \tilde{\psi}_0(\theta)|^2 d\theta \int_{S^2} |\nabla_n (\mathbf{p} \cdot \mathbf{n})|^2 d\mathbf{n}. \end{aligned}$$

Therefore by density of smooth functions in $H_0^1(\text{SO}(3))$, we get that $\bar{\psi} \in H_0^1(\text{SO}(3))$ if and only if $\tilde{\psi}_0 \in H$, where

$$H := \left\{ \psi \left| \int_{(0,\pi)} \psi^2 \sin^2 \theta d\theta < \infty, \int_{(0,\pi)} |\partial_{\theta} (\sin \theta \psi(\theta))|^2 \sin^2(\theta/2) d\theta < \infty \right. \right\}.$$

This Hilbert space is equipped with the corresponding norm:

$$\|\psi\|_H^2 = \int_{(0,\pi)} \psi^2 \sin^2 \theta d\theta + \int_{(0,\pi)} |\partial_\theta(\sin \theta \psi(\theta))|^2 \sin^2(\theta/2) d\theta.$$

Now, Eq. (4.16) written in weak form in H and tested against any $\phi \in H$ reads

$$\begin{aligned} a(\tilde{\psi}_0, \phi) &:= - \int_{(0,\pi)} m(\theta) \left[\sin^2(\theta/2) \partial_\theta(\sin \theta \tilde{\psi}_0(\theta)) \partial_\theta(\sin \theta \phi(\theta)) \right. \\ &\quad \left. + \frac{1}{2} \sin^2 \theta \tilde{\psi}_0(\theta) \phi(\theta) \right] d\theta \\ &= \int_{(0,\pi)} \sin^2 \theta \sin^2(\theta/2) m(\theta) \phi d\theta =: b(\phi). \end{aligned}$$

It holds for some $c, c', c'' > 0$ that $|a(\psi, \phi)| \leq c \|\psi\|_H \|\phi\|_H$ since $m = m(\theta)$ is bounded; and also $|a(\psi, \psi)| \geq c' \|\psi\|_H^2$ since there exists $m_0 > 0$ such that $m(\theta) > m_0$ for all $\theta \in [0, \pi]$; finally, we also have that $|b(\phi)| \leq c'' \|\phi\|_H^2$. Therefore, by the Lax–Milgram theorem, there exists a (unique) solution $\tilde{\psi}_0 \in H$ to (4.16), which corresponds to a (unique) $\tilde{\psi}$ in $H_0^1(\text{SO}(3))$. \square

4.5. The macroscopic limit

In this section, we investigate the hydrodynamic limit. To state the theorem we first give the definitions of the first-order operators δ_x and \mathbf{r}_x . For a smooth function Λ from \mathbb{R}^3 to $\text{SO}(3)$, and for $x \in \mathbb{R}^3$, we define the following matrix $\mathcal{D}_x(\Lambda)$ such that for any $\mathbf{w} \in \mathbb{R}^3$, we have

$$(\mathbf{w} \cdot \nabla_x) \Lambda = [\mathcal{D}_x(\Lambda) \mathbf{w}]_\times \Lambda. \quad (4.19)$$

Notice that this first-order differential equation \mathcal{D}_x is well defined as a matrix; for a given vector \mathbf{w} , the matrix $(\mathbf{w} \cdot \nabla_x) \Lambda$ is in T_Λ and thanks to Proposition A.3, it is of the form $P\Lambda$, with P an antisymmetric matrix. Therefore there exists a vector $\mathcal{D}_x(\Lambda)(\mathbf{w}) \in \mathbb{R}^3$ depending on \mathbf{w} such that $P = [\mathcal{D}_x(\Lambda)(\mathbf{w})]_\times$. The function $\mathbf{w} \mapsto \mathcal{D}_x(\Lambda)(\mathbf{w})$ is linear from \mathbb{R}^3 to \mathbb{R}^3 , so $\mathcal{D}_x(\Lambda)$ can be identified as a matrix.

We now define the first-order operators δ_x (scalar) and \mathbf{r}_x (vector), by

$$\delta_x(\Lambda) = \text{tr}(\mathcal{D}_x(\Lambda)) \quad \text{and} \quad [\mathbf{r}_x(\Lambda)]_\times = \mathcal{D}_x(\Lambda) - \mathcal{D}_x(\Lambda)^T. \quad (4.20)$$

We first give an invariance property which allows for a simple expression for these operators.

Proposition 4.7. *The operators \mathcal{D}_x, δ_x and \mathbf{r}_x are right invariant in the following sense: if A is a fixed matrix in $\text{SO}(3)$ and $\Lambda : \mathbb{R}^3 \rightarrow \text{SO}(3)$ a smooth function, we have*

$$\mathcal{D}_x(\Lambda A) = \mathcal{D}_x(\Lambda), \quad \delta_x(\Lambda A) = \delta_x(\Lambda) \quad \text{and} \quad \mathbf{r}_x(\Lambda A) = \mathbf{r}_x(\Lambda).$$

Consequently, in the neighbourhood of $x_0 \in \mathbb{R}^3$, we can write

$$\Lambda(x) = \exp([\mathbf{b}(x)]_\times) \Lambda(x_0)$$

where \mathbf{b} is a smooth function from a neighbourhood of x_0 into \mathbb{R}^3 such that $\mathbf{b}(x_0) = 0$, and we have

$$(\mathcal{D}_x(\Lambda))_{ij}(x_0) = \partial_j \mathbf{b}_i(x_0),$$

and therefore

$$\delta_x(\Lambda)(x_0) = (\nabla_x \cdot \mathbf{b})(x_0) \quad \text{and} \quad \mathbf{r}_\Lambda(x_0) = (\nabla_x \times \mathbf{b})(x_0).$$

Proof. For any $\mathbf{w} \in \mathbb{R}^3$, we have, since A is constant:

$$[\mathcal{D}_x(\Lambda A) \mathbf{w}]_\times \Lambda A = \mathbf{w} \cdot \nabla_x (\Lambda A) = (\mathbf{w} \cdot \nabla_x \Lambda) A = [\mathcal{D}_x(\Lambda) \mathbf{w}]_\times \Lambda A.$$

This proves that $\mathcal{D}_x(\Lambda A) = \mathcal{D}_x(\Lambda)$, and by (4.20), the same is obviously true for δ_x and \mathbf{r}_x .

We now write, in the neighbourhood of x_0 , that $\Lambda(x) = \exp([\mathbf{b}(x)]_\times) \Lambda(x_0)$, with \mathbf{b} smooth in the neighbourhood of x_0 and $\mathbf{b}(x_0) = 0$. Then we have $\mathcal{D}_x(\Lambda) = \mathcal{D}_x(\exp([\mathbf{b}]_\times))$. We perform a Taylor expansion around x_0 of $\exp([\mathbf{b}]_\times)$:

$$\exp([\mathbf{b}(x)]_\times) = \text{Id} + [\mathbf{b}(x)]_\times + M(x),$$

where $M(x)$ is of order 2 in the coordinates $\mathbf{b}_1, \mathbf{b}_2, \mathbf{b}_3$, (since \mathbf{b} is smooth in the neighbourhood of x_0 and $\mathbf{b}(x_0) = 0$), therefore

$$\partial_1 M(x_0) = \partial_2 M(x_0) = \partial_3 M(x_0) = 0.$$

We then get, since $\exp([\mathbf{b}(x_0)]_\times) = \text{Id}$, that

$$[\mathcal{D}_x(\exp([\mathbf{b}]_\times))(x_0) \mathbf{w}]_\times = \mathbf{w} \cdot \nabla_x (\exp([\mathbf{b}]_\times))(x_0) = [(\mathbf{w} \cdot \nabla_x \mathbf{b})(x_0)]_\times,$$

and therefore $\mathcal{D}_x(\Lambda)(x_0) \mathbf{w} = \mathcal{D}_x(\exp([\mathbf{b}]_\times))(x_0) \mathbf{w} = (\mathbf{w} \cdot \nabla_x \mathbf{b})(x_0)$. Taking $\mathbf{w} = \mathbf{e}_j$, we get $\mathcal{D}_x(\Lambda)(x_0) \mathbf{e}_j = \partial_j \mathbf{b}(x_0)$, and thus $(\mathcal{D}_x(\Lambda)(x_0))_{ij} = \mathbf{e}_i \cdot \mathcal{D}_x(\Lambda)(x_0) \mathbf{e}_j = \partial_j \mathbf{b}_i$. The formula for $\delta_x(\Lambda)$ follows from (4.20), since $\nabla_x \cdot \mathbf{b} = \sum_i \partial_i \mathbf{b}_i$. Finally by the definition of $[\cdot]_\times$ (see (3.2)), we get

$$[\nabla_x \times \mathbf{b}]_\times = \begin{pmatrix} 0 & \partial_2 \mathbf{b}_1 - \partial_1 \mathbf{b}_2 & \partial_3 \mathbf{b}_1 - \partial_1 \mathbf{b}_3 \\ \partial_1 \mathbf{b}_2 - \partial_2 \mathbf{b}_1 & 0 & \partial_3 \mathbf{b}_2 - \partial_2 \mathbf{b}_3 \\ \partial_1 \mathbf{b}_3 - \partial_3 \mathbf{b}_1 & \partial_2 \mathbf{b}_3 - \partial_3 \mathbf{b}_2 & 0 \end{pmatrix},$$

so from (4.20) we obtain $(\nabla_x \times \mathbf{b})(x_0) = \mathbf{r}_x(\Lambda)(x_0)$. □

We are now ready to state the main theorem of our paper (see Sec. 2 for a discussion on this result).

Theorem 4.1. ((Formal) macroscopic limit) *When $\varepsilon \rightarrow 0$ in the kinetic equation (4.1) it holds (formally) that*

$$f_\varepsilon \rightarrow f = f(x, A, t) = \rho M_\Lambda(A), \quad \Lambda = \Lambda(t, x) \in \text{SO}(3), \quad \rho = \rho(t, x) \geq 0.$$

Moreover, if this convergence is strong enough and the functions Λ and ρ are smooth enough, they satisfy the following first-order system of partial differential equations:

$$\partial_t \rho + \nabla_x \cdot (c_1 \rho \Lambda \mathbf{e}_1) = 0, \quad (4.21)$$

$$\begin{aligned} \rho(\partial_t \Lambda + c_2((\Lambda \mathbf{e}_1) \cdot \nabla_x) \Lambda) + [(\Lambda \mathbf{e}_1) \times (c_3 \nabla_x \rho + c_4 \rho \mathbf{r}_x(\Lambda)) \\ + c_4 \rho \delta_x(\Lambda) \Lambda \mathbf{e}_1] \times \Lambda = 0, \end{aligned} \quad (4.22)$$

where $c_1 = c_1(\nu, d) = \frac{2}{3} \langle \frac{1}{2} + \cos \theta \rangle_{m(\theta) \sin^2(\theta/2)}$ is the constant given in (4.7) and

$$\begin{aligned} c_2 &= \frac{1}{5} \langle 2 + 3 \cos \theta \rangle_{\tilde{m}(\theta) \sin^2(\theta/2)}, \\ c_3 &= d \left\langle \nu \left(\frac{1}{2} + \cos \theta \right)^{-1} \right\rangle_{\tilde{m}(\theta) \sin^2(\theta/2)}, \\ c_4 &= \frac{1}{5} \langle 1 - \cos \theta \rangle_{\tilde{m}(\theta) \sin^2(\theta/2)}, \end{aligned}$$

where the notation $\langle \cdot \rangle_{\tilde{m}(\theta) \sin^2(\theta/2)}$ is defined in (4.6). The function

$$\tilde{m} : (0, \pi) \rightarrow (0, +\infty)$$

is given by

$$\tilde{m}(\theta) := \nu \left(\frac{1}{2} + \cos \theta \right) \sin^2 \theta m(\theta) \tilde{\psi}_0(\theta), \quad (4.23)$$

where $m(\theta) = \exp(d^{-1} \sigma(\frac{1}{2} + \cos \theta))$ is the same as in (4.8) and $\tilde{\psi}_0$ is the solution of Eq. (4.16).

Proof. Suppose that $f_\varepsilon \rightarrow f$ as $\varepsilon \rightarrow 0$, then using (4.1) we get $Q(f_\varepsilon) = \mathcal{O}(\varepsilon)$, which formally yields $Q(f) = 0$ and by Lemma 4.3 we have that

$$f = f(x, A, t) = \rho M_\Lambda(A), \quad \text{with} \quad \Lambda = \Lambda(t, x) \in \text{SO}(3), \quad \rho = \rho(t, x) \geq 0.$$

Using the conservation of mass (integrating (4.1) on $\text{SO}(3)$), we have that

$$\partial_t \rho_\varepsilon + \nabla_x \cdot j[f_\varepsilon] = \mathcal{O}(\varepsilon),$$

where

$$\rho_\varepsilon(t, x) := \int_{\text{SO}(3)} f_\varepsilon(x, A, t) dA, \quad j[f_\varepsilon] := \int_{\text{SO}(3)} A \mathbf{e}_1 f_\varepsilon dA,$$

and in the limit (formally)

$$\rho_\varepsilon \rightarrow \rho,$$

$$j[f_\varepsilon] \rightarrow \rho \int_{\text{SO}(3)} A \mathbf{e}_1 M_\Lambda(A) dA = \rho \lambda[M_\Lambda] \mathbf{e}_1 = \rho c_1 \Lambda \mathbf{e}_1,$$

thanks to Lemma 4.4. This gives us the continuity equation (4.21) for ρ .

Now, we want to obtain the equation for Λ . We write $\Lambda^\varepsilon = \Lambda[f^\varepsilon]$, and we take $P \in \mathcal{A}$ a given antisymmetric matrix. We consider the non-constant GCI

associated to Λ^ε and corresponding to P in (4.18): $\psi^\varepsilon(A) = P \cdot ((\Lambda^\varepsilon)^T A) \bar{\psi}_0(\Lambda^\varepsilon \cdot A)$. Since we have $\psi^\varepsilon \in \text{GCI}(\Lambda[f^\varepsilon])$, we obtain, thanks to the main property (4.9) of the GCI, that

$$\int_{\text{SO}(3)} Q(f^\varepsilon) \psi^\varepsilon dA = 0.$$

Multiplying (4.1) by ψ^ε , integrating with respect to A on $\text{SO}(3)$ and using the expression of ψ^ε as stated above, we obtain

$$\int_{\text{SO}(3)} (\partial_t f^\varepsilon + A \mathbf{e}_1 \cdot \nabla_x f^\varepsilon + \mathcal{O}(\varepsilon)) P \cdot ((\Lambda^\varepsilon)^T A) \bar{\psi}_0(\Lambda^\varepsilon \cdot A) dA = 0.$$

Assuming the convergence $f^\varepsilon \rightarrow f$ is sufficiently strong, we get in the limit

$$\int_{\text{SO}(3)} (\partial_t(\rho M_\Lambda) + A \mathbf{e}_1 \cdot \nabla_x(\rho M_\Lambda))(P \cdot \Lambda^T A) \bar{\psi}_0(\Lambda \cdot A) dA = 0. \quad (4.24)$$

Since (4.24) is true for any $P \in \mathcal{A}$, the matrix

$$\int_{\text{SO}(3)} (\partial_t(\rho M_\Lambda) + A \mathbf{e}_1 \cdot \nabla_x(\rho M_\Lambda)) \bar{\psi}_0(\Lambda \cdot A) \Lambda^T A dA = 0$$

is orthogonal to all antisymmetric matrices. Therefore, it must be a symmetric matrix, meaning that we have

$$X := \int_{\text{SO}(3)} (\partial_t(\rho M_\Lambda) + A \mathbf{e}_1 \cdot \nabla_x(\rho M_\Lambda)) \bar{\psi}_0(\Lambda \cdot A) (\Lambda^T A - A^T \Lambda) dA = 0. \quad (4.25)$$

We have with the definition of M_Λ in (4.4) that:

$$\partial_t(\rho M_\Lambda) = M_\Lambda(\partial_t \rho + d^{-1} \nu(\Lambda \cdot A) \rho(A \cdot \partial_t \Lambda)),$$

$$(A \mathbf{e}_1 \cdot \nabla_x)(\rho M_\Lambda) = M_\Lambda(A \mathbf{e}_1 \cdot \nabla_x \rho + d^{-1} \nu(\Lambda \cdot A) \rho(A \cdot (A \mathbf{e}_1 \cdot \nabla_x) \Lambda)).$$

Inserting the two previous expressions into (4.25), we compute separately each component of X defined by:

$$X_1 := \int_{\text{SO}(3)} \partial_t \rho M_\Lambda \bar{\psi}_0(\Lambda \cdot A) (\Lambda^T A - A^T \Lambda) dA,$$

$$X_2 := \int_{\text{SO}(3)} d^{-1} \nu(\Lambda \cdot A) \rho(A \cdot \partial_t \Lambda) M_\Lambda \bar{\psi}_0(\Lambda \cdot A) (\Lambda^T A - A^T \Lambda) dA,$$

$$X_3 := \int_{\text{SO}(3)} A \mathbf{e}_1 \cdot \nabla_x \rho M_\Lambda \bar{\psi}_0(\Lambda \cdot A) (\Lambda^T A - A^T \Lambda) dA,$$

$$X_4 := \int_{\text{SO}(3)} d^{-1} \nu(\Lambda \cdot A) \rho(A \cdot (A \mathbf{e}_1 \cdot \nabla_x) \Lambda) M_\Lambda \bar{\psi}_0(\Lambda \cdot A) (\Lambda^T A - A^T \Lambda) dA,$$

so $X = X_1 + X_2 + X_3 + X_4$.

For the first term we have (changing variables $B = \Lambda^T A$):

$$X_1 = \partial_t \rho \int_{\text{SO}(3)} M_{\text{Id}}(B) \bar{\psi}_0(\text{Id} \cdot B) (B - B^T) dB = 0,$$

since both $M_{\text{Id}}(B)$ and $\bar{\psi}_0(\text{Id} \cdot B)$ are invariant by the change $B \mapsto B^T$.

For the term X_2 we make the change of variables $B = \Lambda^T A$ and compute

$$\begin{aligned} X_2 &= \rho \int_{\text{SO}(3)} d^{-1} \nu(\text{Id} \cdot B) (\Lambda B \cdot \partial_t \Lambda) M_{\text{Id}}(B) \bar{\psi}_0(\text{Id} \cdot B) (B - B^T) dB \\ &= \frac{2d^{-1} \rho}{\pi Z} \int_{(0, \pi) \times S^2} (\Lambda(\text{Id} + \sin \theta [\mathbf{n}]_{\times} + (1 - \cos \theta) [\mathbf{n}]_{\times}^2)) \cdot \partial_t \Lambda \\ &\quad \times \sin^2(\theta/2) \nu \left(\frac{1}{2} + \cos \theta \right) m(\theta) \tilde{\psi}_0(\theta) 2 \sin \theta [\mathbf{n}]_{\times} d\theta d\mathbf{n}, \end{aligned}$$

where we have used the expression of the Haar measure $dB = \frac{2}{\pi} \sin^2(\theta/2) d\theta d\mathbf{n}$ (see Lemma 4.2) and that writing $B = B(\theta, \mathbf{n}) = \text{Id} + \sin \theta [\mathbf{n}]_{\times} + (1 - \cos \theta) [\mathbf{n}]_{\times}^2$ thanks to Rodrigues' formula (3.3), we have $B - B^T = 2 \sin \theta [\mathbf{n}]_{\times}$. Removing odd terms with respect to the change $\mathbf{n} \mapsto -\mathbf{n}$, we obtain

$$\begin{aligned} X_2 &= \frac{4d^{-1} \rho}{\pi Z} \int_{(0, \pi) \times S^2} \nu \left(\frac{1}{2} + \cos \theta \right) \sin^2 \theta m(\theta) \tilde{\psi}_0(\theta) \\ &\quad \times \sin^2(\theta/2) (\Lambda [\mathbf{n}]_{\times} \cdot \partial_t \Lambda) [\mathbf{n}]_{\times} d\theta d\mathbf{n}. \end{aligned}$$

Now since $\partial_t \Lambda \in T_{\Lambda}$, we have $\Lambda^T \partial_t \Lambda \in \mathcal{A}$ (antisymmetric, see Proposition A.2), and so

$$\Lambda^T \partial_t \Lambda = [\boldsymbol{\lambda}_t]_{\times},$$

for some vector $\boldsymbol{\lambda}_t$. Therefore

$$(\Lambda [\mathbf{n}]_{\times}) \cdot \partial_t \Lambda = [\mathbf{n}]_{\times} \cdot (\Lambda^T \partial_t \Lambda) = [\mathbf{n}]_{\times} \cdot [\boldsymbol{\lambda}_t]_{\times} = (\mathbf{n} \cdot \boldsymbol{\lambda}_t).$$

So using the definition (4.23) of $\tilde{m}(\theta)$, we get

$$\begin{aligned} X_2 &= \frac{4d^{-1} \rho}{\pi Z} \int_{(0, \pi) \times S^2} \tilde{m}(\theta) \sin^2(\theta/2) (\mathbf{n} \cdot \boldsymbol{\lambda}_t) [\mathbf{n}]_{\times} d\theta d\mathbf{n} \\ &= \frac{4d^{-1} \rho}{\pi Z} \left[\int_{(0, \pi) \times S^2} \tilde{m}(\theta) \sin^2(\theta/2) (\mathbf{n} \cdot \boldsymbol{\lambda}_t) \mathbf{n} d\theta d\mathbf{n} \right]_{\times} \\ &= \frac{4d^{-1} \rho}{3\pi Z} \left(\int_0^{\pi} \tilde{m}(\theta) \sin^2(\theta/2) d\theta \right) [\boldsymbol{\lambda}_t]_{\times}, \end{aligned}$$

because the mapping $\mathbf{w} \mapsto [\mathbf{w}]_{\times}$ is linear, and $\int_{S^2} \mathbf{n} \otimes \mathbf{n} d\mathbf{n} = \frac{1}{3} \text{Id}$.

Denote by

$$C_2 := \frac{4d^{-1}}{3\pi Z} \left(\int_0^{\pi} \tilde{m}(\theta) \sin^2(\theta/2) d\theta \right),$$

then we conclude that

$$X_2 = C_2 \rho \Lambda^T \partial_t \Lambda.$$

Now, for the term X_3 we compute the following, starting again by the change of variables $B = \Lambda^T A$:

$$\begin{aligned} X_3 &= \int_{\text{SO}(3)} (\Lambda B \mathbf{e}_1 \cdot \nabla_x \rho) M_{\text{Id}}(B) \bar{\psi}_0(\text{Id} \cdot B) (B - B^T) dB \\ &= \frac{4}{\pi Z} \int_{(0, \pi) \times S^2} m(\theta) \tilde{\psi}_0(\theta) \sin \theta \sin^2(\theta/2) \\ &\quad \times (\Lambda(\text{Id} + \sin \theta [\mathbf{n}]_{\times} + (1 - \cos \theta) [\mathbf{n}]_{\times}^2) \mathbf{e}_1 \cdot \nabla_x \rho) [\mathbf{n}]_{\times} d\theta d\mathbf{n} \\ &= \frac{4}{\pi Z} \int_{(0, \pi) \times S^2} m(\theta) \tilde{\psi}_0(\theta) \sin^2 \theta \sin^2(\theta/2) (\Lambda [\mathbf{n}]_{\times} \mathbf{e}_1 \cdot \nabla_x \rho) [\mathbf{n}]_{\times} d\theta d\mathbf{n} \\ &= \frac{4}{\pi Z} \left[\int_{(0, \pi) \times S^2} \frac{\tilde{m}(\theta)}{\nu(\frac{1}{2} + \cos \theta)} \sin^2(\theta/2) (\mathbf{n} \cdot (\mathbf{e}_1 \times \Lambda^T \nabla_x \rho)) \mathbf{n} d\theta d\mathbf{n} \right]_{\times} \\ &= \frac{4}{3\pi Z} \left(\int_0^{\pi} \frac{\tilde{m}(\theta)}{\nu(\frac{1}{2} + \cos \theta)} \sin^2(\theta/2) d\theta \right) [\mathbf{e}_1 \times \Lambda^T \nabla_x \rho]_{\times}, \end{aligned}$$

where we used similar considerations as for X_2 , as well as that

$$\Lambda [\mathbf{n}]_{\times} \mathbf{e}_1 \cdot \nabla_x \rho = [\mathbf{n}]_{\times} \mathbf{e}_1 \cdot (\Lambda^T \nabla_x \rho) = (\mathbf{n} \times \mathbf{e}_1) \cdot (\Lambda^T \nabla_x \rho) = \mathbf{n} \cdot (\mathbf{e}_1 \times \Lambda^T \nabla_x \rho).$$

Denote by

$$C_3 := \frac{4}{3\pi Z} \left(\int_0^{\pi} \frac{\tilde{m}(\theta)}{\nu(\frac{1}{2} + \cos \theta)} \sin^2(\theta/2) d\theta \right),$$

then

$$X_3 = C_3 [\mathbf{e}_1 \times \Lambda^T \nabla_x \rho]_x.$$

We now compute X_4 in the same way, with the change of variables $B = \Lambda^T A$:

$$X_4 = \rho d^{-1} \int_{\text{SO}(3)} (\nu(\text{Id} \cdot B) (\Lambda B \cdot (\Lambda B \mathbf{e}_1 \cdot \nabla_x) \Lambda)) M_{\text{Id}}(B) \bar{\psi}_0(\text{Id} \cdot B) (B - B^T) dB.$$

We now use the definition of $\mathcal{D}_x(\Lambda)$ given in (4.19) to get

$$\begin{aligned} X_4 &= \rho d^{-1} \int_{\text{SO}(3)} (\nu(\text{Id} \cdot B) (\Lambda B \cdot ([\mathcal{D}_x(\Lambda) \Lambda B \mathbf{e}_1]_{\times} \Lambda))) \\ &\quad \times M_{\text{Id}}(B) (B - B^T) \bar{\psi}_0(\text{Id} \cdot B) dB. \end{aligned}$$

Using the fact that $\Lambda^T [\mathbf{w}]_{\times} = [\Lambda^T \mathbf{w}]_{\times} \Lambda^T$ for all $\mathbf{w} \in \mathbb{R}^3$, we have

$$\Lambda B \cdot ([\mathcal{D}_x(\Lambda) \Lambda B \mathbf{e}_1]_{\times} \Lambda) = B \cdot [\Lambda^T \mathcal{D}_x(\Lambda) \Lambda B \mathbf{e}_1]_{\times}.$$

To simplify the notations, we denote $L = \Lambda^T \mathcal{D}_x(\Lambda) \Lambda$. Since the symmetric part of B does not contribute to the scalar product $B \cdot [L B \mathbf{e}_1]_\times$, we get

$$\Lambda B \cdot ([\mathcal{D}_x(\Lambda) \Lambda B \mathbf{e}_1]_\times \Lambda) = B \cdot [L B \mathbf{e}_1]_\times = \sin \theta [\mathbf{n}]_\times \cdot [L B \mathbf{e}_1]_\times = \sin \theta \mathbf{n} \cdot L B \mathbf{e}_1.$$

Therefore we obtain, in the same manner as before,

$$\begin{aligned} X_4 &= \frac{4\rho d^{-1}}{\pi Z} \int_0^\pi \tilde{m}(\theta) \sin^2(\theta/2) \\ &\quad \times \left[\int_{S^2} (\mathbf{n} \cdot (L(\text{Id} + \sin \theta [\mathbf{n}]_\times + (1 - \cos \theta) [\mathbf{n}]_\times^2 \mathbf{e}_1)) \mathbf{n} d\mathbf{n} \right]_\times d\theta, \end{aligned}$$

and we have to know the value of

$$\begin{aligned} \mathbf{y}(\theta) &:= \int_{S^2} (\mathbf{n} \cdot (L(\text{Id} + \sin \theta [\mathbf{n}]_\times + (1 - \cos \theta) [\mathbf{n}]_\times^2 \mathbf{e}_1)) \mathbf{n} d\mathbf{n} \\ &= \int_{S^2} (\mathbf{n} \cdot (L(\cos \theta \mathbf{e}_1 + (1 - \cos \theta)(\mathbf{n} \cdot \mathbf{e}_1) \mathbf{n}))) \mathbf{n} d\mathbf{n} \\ &= \frac{1}{3} \cos \theta L \mathbf{e}_1 + (1 - \cos \theta) \left(\int_{S^2} \mathbf{n} \cdot L \mathbf{n} (\mathbf{n} \otimes \mathbf{n}) d\mathbf{n} \right) \mathbf{e}_1, \end{aligned}$$

where the term involving $[\mathbf{n}]_\times$ vanishes since its integrand is odd with respect to $\mathbf{n} \mapsto -\mathbf{n}$.

To compute the second term of this expression we will make use of the following lemma proved at the end of this section.

Lemma 4.5. *For a given matrix $L \in \mathcal{M}$, we have*

$$\int_{S^2} \mathbf{n} \cdot L \mathbf{n} (\mathbf{n} \otimes \mathbf{n}) d\mathbf{n} = \frac{1}{15} (L + L^T) + \frac{1}{15} \text{tr}(L) \text{Id}.$$

Using this lemma we have that

$$\begin{aligned} \mathbf{y}(\theta) &= \frac{1}{3} \cos \theta L \mathbf{e}_1 + (1 - \cos \theta) \left(\frac{1}{15} (L + L^T) + \frac{1}{15} \text{tr}(L) \text{Id} \right) \mathbf{e}_1 \\ &= \frac{1}{15} (1 + 4 \cos \theta) L \mathbf{e}_1 + \frac{1}{15} (1 - \cos \theta) (L^T \mathbf{e}_1 + \text{tr}(L) \mathbf{e}_1). \end{aligned}$$

Therefore we obtain

$$\begin{aligned} X_4 &= \frac{4\rho d^{-1}}{\pi Z} \int_0^\pi \tilde{m}(\theta) \sin^2(\theta/2) [\mathbf{y}(\theta)]_\times d\theta \\ &= \frac{4\rho d^{-1}}{15\pi Z} \int_0^\pi \tilde{m}(\theta) \sin^2(\theta/2) ((1 + 4 \cos \theta) [L \mathbf{e}_1]_\times \\ &\quad + (1 - \cos \theta) [L^T \mathbf{e}_1 + \text{tr}(L) \mathbf{e}_1]_\times) d\theta \\ &= \rho (C_4 [L \mathbf{e}_1]_\times + C_5 [L^T \mathbf{e}_1 + \text{tr}(L) \mathbf{e}_1]_\times), \end{aligned}$$

for

$$C_4 := \frac{4d^{-1}}{15\pi Z} \int_0^\pi \tilde{m}(\theta) \sin^2(\theta/2)(1 + 4 \cos \theta) d\theta,$$

$$C_5 := \frac{4d^{-1}}{15\pi Z} \int_0^\pi \tilde{m}(\theta) \sin^2(\theta/2)(1 - \cos \theta) d\theta.$$

Finally putting all the terms together we have that

$$\begin{aligned} 0 &= X = X_1 + X_2 + X_3 + X_4 \\ &= C_2 \rho \Lambda^T \partial_t \Lambda + C_3 [\mathbf{e}_1 \times \Lambda^T \nabla_x \rho]_\times + \rho C_4 [L \mathbf{e}_1]_\times \\ &\quad + \rho C_5 [L^T \mathbf{e}_1 + \text{tr}(L) \mathbf{e}_1]_\times. \end{aligned}$$

In particular $\Lambda X = 0$ and from the fact that $\Lambda[\mathbf{w}]_\times = [\Lambda \mathbf{w}]_\times \Lambda$ for any $\mathbf{w} \in \mathbb{R}^3$ we get

$$\begin{aligned} 0 = \Lambda X &= C_2 \rho \partial_t \Lambda + C_3 [(\Lambda \mathbf{e}_1) \times \nabla_x \rho]_\times \Lambda + C_4 \rho [\Lambda L \mathbf{e}_1]_\times \Lambda \\ &\quad + C_5 \rho [\Lambda L^T \mathbf{e}_1 + \text{tr}(L) \Lambda \mathbf{e}_1]_\times \Lambda. \end{aligned} \quad (4.26)$$

Since we have taken $L = \Lambda^T \mathcal{D}_x(\Lambda) \Lambda$, we get that $\text{tr}(L) = \text{tr}(\mathcal{D}_x(\Lambda)) = \delta_x(\Lambda)$ and, thanks to (4.20):

$$[\Lambda L^T \mathbf{e}_1]_\times = [\mathcal{D}_x(\Lambda)^T \Lambda \mathbf{e}_1]_\times = [(\mathcal{D}_x(\Lambda) - [\mathbf{r}_x(\Lambda)]_\times) \Lambda \mathbf{e}_1]_\times.$$

Furthermore, we have $[\Lambda L \mathbf{e}_1]_\times \Lambda = [\mathcal{D}_x(\Lambda) \Lambda \mathbf{e}_1]_\times \Lambda = ((\Lambda \mathbf{e}_1) \cdot \nabla_x) \Lambda$ thanks to the definition of \mathcal{D}_x given in (4.19). Finally, inserting these expressions into (4.26) and dividing by C_2 , we get the equation

$$\begin{aligned} &\rho(\partial_t \Lambda + c_2((\Lambda \mathbf{e}_1) \cdot \nabla_x) \Lambda) + c_3[(\Lambda \mathbf{e}_1) \times \nabla_x \rho]_\times \Lambda \\ &\quad + c_4 \rho[-\mathbf{r}_x(\Lambda) \times (\Lambda \mathbf{e}_1) + \delta_x(\Lambda) \Lambda \mathbf{e}_1]_\times \Lambda = 0, \end{aligned}$$

for

$$\begin{aligned} c_2 &= \frac{C_4 + C_5}{C_2} = \frac{1}{5} \langle 2 + 3 \cos \theta \rangle_{\tilde{m}(\theta) \sin^2(\theta/2)}, \\ c_3 &= \frac{C_3}{C_2} = d \left\langle \nu \left(\frac{1}{2} + \cos \theta \right)^{-1} \right\rangle_{\tilde{m}(\theta) \sin^2(\theta/2)}, \\ c_4 &= \frac{C_5}{C_2} = \frac{1}{5} \langle 1 - \cos \theta \rangle_{\tilde{m}(\theta) \sin^2(\theta/2)}, \end{aligned}$$

which ends the proof. \square

Proof of Lemma 4.5. Denote by $\mathcal{I}(L)$ the integral that we want to compute

$$\mathcal{I}(L) := \int_{S^2} \mathbf{n} \cdot L \mathbf{n} (\mathbf{n} \otimes \mathbf{n}) d\mathbf{n},$$

then, written in components, we have

$$\begin{aligned}\mathcal{I}(L)_{ij} &= \int_{S^2} (\mathbf{n} \cdot L\mathbf{n})(\mathbf{e}_i \cdot \mathbf{n})(\mathbf{e}_j \cdot \mathbf{n})d\mathbf{n} \\ &= \begin{cases} (L_{ij} + L_{ji}) \int_{S^2} (\mathbf{e}_i \cdot \mathbf{n})^2 (\mathbf{e}_j \cdot \mathbf{n})^2 d\mathbf{n} & \text{if } i \neq j, \\ \sum_k L_{kk} \int_{S^2} (\mathbf{e}_k \cdot \mathbf{n})^2 (\mathbf{e}_i \cdot \mathbf{n})^2 d\mathbf{n} & \text{if } i = j \end{cases} \\ &= \begin{cases} \frac{1}{15}(L_{ij} + L_{ji}) & \text{if } i \neq j, \\ \frac{1}{15} \sum_k L_{kk} + \frac{2}{15}L_{ii} & \text{if } i = j \end{cases} \\ &= \frac{1}{15}(L_{ij} + L_{ji}) + \begin{cases} 0 & \text{if } i \neq j, \\ \frac{1}{15} \sum_k L_{kk} & \text{if } i = j, \end{cases}\end{aligned}$$

from which we conclude the lemma. In the computations we used that:

$$\begin{aligned}\text{for } i \neq j, \quad \int_{S^2} (\mathbf{e}_i \cdot \mathbf{n})^2 (\mathbf{e}_j \cdot \mathbf{n})^2 d\mathbf{n} &= \frac{1}{4\pi} \int_{[0,\pi] \times [0,2\pi]} \sin^3 \phi \cos^2 \psi \cos^2 \phi d\phi d\psi = \frac{1}{15}; \\ \text{for } k = i, \quad \int_{S^2} (\mathbf{e}_k \cdot \mathbf{n})^4 d\mathbf{n} &= \frac{1}{4\pi} \int_{[0,\pi] \times [0,2\pi]} \cos^4 \phi \sin \phi d\phi d\psi = \frac{1}{5}. \quad \square\end{aligned}$$

Finally, we consider the orthonormal basis given by

$$\{\Lambda \mathbf{e}_1 =: \Omega, \Lambda \mathbf{e}_2 =: \mathbf{u}, \Lambda \mathbf{e}_3 =: \mathbf{v}\},$$

where $\{\mathbf{e}_1, \mathbf{e}_2, \mathbf{e}_3\}$ is the canonical basis of \mathbb{R}^3 . We can have an expression of the operators δ_x and \mathbf{r}_x in terms of these unit vectors $\{\Omega, \mathbf{u}, \mathbf{v}\}$, which allows to rewrite the evolution equation of Λ as three evolution equations for these vectors.

Proposition 4.8. *We have:*

$$\delta_x(\Lambda) = [(\Omega \cdot \nabla_x)\mathbf{u}] \cdot \mathbf{v} + [(\mathbf{u} \cdot \nabla_x)\mathbf{v}] \cdot \Omega + [(\mathbf{v} \cdot \nabla_x)\Omega] \cdot \mathbf{u}, \quad (4.27)$$

$$\mathbf{r}_x(\Lambda) = (\nabla_x \cdot \Omega)\Omega + (\nabla_x \cdot \mathbf{u})\mathbf{u} + (\nabla_x \cdot \mathbf{v})\mathbf{v}. \quad (4.28)$$

Consequently, we have the following evolution equations for Ω, \mathbf{u} and \mathbf{v} , corresponding to the evolution equation of Λ given in (4.22):

$$\rho D_t \Omega + P_{\Omega^\perp} (c_3 \nabla_x \rho + c_4 \rho ((\nabla_x \cdot \mathbf{u})\mathbf{u} + (\nabla_x \cdot \mathbf{v})\mathbf{v})) = 0, \quad (4.29)$$

$$\rho D_t \mathbf{u} - (c_3 \mathbf{u} \cdot \nabla_x \rho + c_4 \rho \nabla_x \cdot \mathbf{u})\Omega + c_4 \rho \delta_x(\Omega, \mathbf{u}, \mathbf{v})\mathbf{v} = 0, \quad (4.30)$$

$$\rho D_t \mathbf{v} - (c_3 \mathbf{v} \cdot \nabla_x \rho + c_4 \rho \nabla_x \cdot \mathbf{v})\Omega - c_4 \rho \delta_x(\Omega, \mathbf{u}, \mathbf{v})\mathbf{u} = 0, \quad (4.31)$$

where $D_t := \partial_t + c_2(\Omega \cdot \nabla_x)$, and where $\delta_x(\Omega, \mathbf{u}, \mathbf{v})$ is the expression of $\delta_x(\Lambda)$ given by (4.27).

Proof. We first prove (4.27). We have

$$\begin{aligned}
 \delta_x(\Lambda) &= \text{tr}(\mathcal{D}_x(\Lambda)) = \text{tr}(\Lambda^T \mathcal{D}_x(\Lambda) \Lambda) \\
 &= \sum_k \Lambda^T \mathcal{D}_x(\Lambda) \Lambda \mathbf{e}_k \cdot \mathbf{e}_k = \sum_k (\mathcal{D}_x(\Lambda) \Lambda \mathbf{e}_k) \cdot \Lambda \mathbf{e}_k \\
 &= \sum_k [\mathcal{D}_x(\Lambda) \Lambda \mathbf{e}_k]_{\times} \cdot [\Lambda \mathbf{e}_k]_{\times} = \sum_k [\mathcal{D}_x(\Lambda) \Lambda \mathbf{e}_k]_{\times} \Lambda \cdot [\Lambda \mathbf{e}_k]_{\times} \Lambda \\
 &= \sum_k ((\Lambda \mathbf{e}_k \cdot \nabla_x) \Lambda) \cdot [\Lambda \mathbf{e}_k]_{\times} \Lambda,
 \end{aligned}$$

thanks to the definition of \mathcal{D}_x given in (4.19). Now we use the fact that for two matrices A, B , we have $A \cdot B = \frac{1}{2} \text{tr}(A^T B) = \frac{1}{2} \sum_i A \mathbf{e}_i \cdot B \mathbf{e}_i$ (half the sum of the scalar products of the corresponding columns of the matrices A and B), to get

$$\begin{aligned}
 \delta_x(\Lambda) &= \frac{1}{2} \sum_k \sum_i [(\Lambda \mathbf{e}_k \cdot \nabla_x) (\Lambda \mathbf{e}_i)] \cdot [(\Lambda \mathbf{e}_k) \times (\Lambda \mathbf{e}_i)] \\
 &= \frac{1}{2} ((\Omega \cdot \nabla_x) \mathbf{u} \cdot \mathbf{v} - (\mathbf{u} \cdot \nabla_x) \Omega \cdot \mathbf{v} - (\Omega \cdot \nabla_x) \mathbf{v} \cdot \mathbf{u} \\
 &\quad + (\mathbf{v} \cdot \nabla_x) \Omega \cdot \mathbf{u} + (\mathbf{u} \cdot \nabla_x) \mathbf{v} \cdot \Omega - (\mathbf{v} \cdot \nabla_x) \mathbf{u} \cdot \Omega) \\
 &= [(\Omega \cdot \nabla_x) \mathbf{u}] \cdot \mathbf{v} + [(\mathbf{u} \cdot \nabla_x) \mathbf{v}] \cdot \Omega + [(\mathbf{v} \cdot \nabla_x) \Omega] \cdot \mathbf{u}.
 \end{aligned}$$

For this last equality we used the fact that

$$0 = (\Omega \cdot \nabla_x)(\mathbf{u} \cdot \mathbf{v}) = (\Omega \cdot \nabla_x) \mathbf{u} \cdot \mathbf{v} + (\Omega \cdot \nabla_x) \mathbf{v} \cdot \mathbf{u}$$

since $\mathbf{u} \perp \mathbf{v}$ and analogously for the other components.

We proceed next to proving the expression of $\mathbf{r}_x(\Lambda)$ given by (4.28). We first prove that $\mathbf{r}_x(\Lambda) \cdot \Omega = \nabla_x \cdot \Omega$. We have (recall that $[\mathbf{r}_x(\Lambda)]_{\times} = \mathcal{D}_x(\Lambda) - \mathcal{D}_x(\Lambda)^T$ and that for all \mathbf{w} in \mathbb{R}^3 , $\mathbf{w} \cdot \nabla_x \Lambda = [\mathcal{D}_x(\Lambda) \mathbf{w}]_{\times} \Lambda$):

$$\begin{aligned}
 \mathbf{r}_x(\Lambda) \cdot \Omega &= \mathbf{r}_x(\Lambda) \cdot (\mathbf{u} \times \mathbf{v}) = \mathbf{v} \cdot ([\mathbf{r}_x(\Lambda)]_{\times} \mathbf{u}) = \mathbf{v} \cdot (\mathcal{D}_x(\Lambda) - \mathcal{D}_x(\Lambda)^T) \mathbf{u} \\
 &= \mathbf{v} \cdot \mathcal{D}_x(\Lambda) \mathbf{u} - \mathbf{u} \cdot \mathcal{D}_x(\Lambda) \mathbf{v} \\
 &= (\Omega \times \mathbf{u}) \cdot \mathcal{D}_x(\Lambda) \mathbf{u} + (\Omega \times \mathbf{v}) \cdot \mathcal{D}_x(\Lambda) \mathbf{u} \\
 &= [\mathcal{D}_x(\Lambda) \mathbf{u}]_{\times} \Omega \cdot \mathbf{u} + [\mathcal{D}_x(\Lambda) \mathbf{v}]_{\times} \Omega \cdot \mathbf{v} \\
 &= [\mathcal{D}_x(\Lambda) \mathbf{u}]_{\times} \Lambda \mathbf{e}_1 \cdot \mathbf{u} + [\mathcal{D}_x(\Lambda) \mathbf{v}]_{\times} \Lambda \mathbf{e}_1 \cdot \mathbf{v} \\
 &= ((\mathbf{u} \cdot \nabla_x) \Lambda \mathbf{e}_1) \cdot \mathbf{u} + ((\mathbf{v} \cdot \nabla_x) \Lambda \mathbf{e}_1) \cdot \mathbf{v} \\
 &= ((\mathbf{u} \cdot \nabla_x) \Omega) \cdot \mathbf{u} + ((\mathbf{v} \cdot \nabla_x) \Omega) \cdot \mathbf{v}.
 \end{aligned}$$

Since $(\Omega \cdot \nabla_x) \Omega$ is orthogonal to Ω , we therefore get

$$\begin{aligned}
 \mathbf{r}_x(\Lambda) \cdot \Omega &= ((\Omega \cdot \nabla_x) \Omega) \cdot \Omega + ((\mathbf{u} \cdot \nabla_x) \Omega) \cdot \mathbf{u} + ((\mathbf{v} \cdot \nabla_x) \Omega) \cdot \mathbf{v} \\
 &= \sum_{i,k,j} \Lambda_{ik} \partial_i \Omega_j \Lambda_{jk} = \sum_{i,j} \partial_i \Omega_j \sum_k \Lambda_{ik} \Lambda_{kj}^T = \sum_i \partial_i \Omega_i = \nabla_x \cdot \Omega,
 \end{aligned}$$

since $\Lambda\Lambda^T = \text{Id}$ (the first line is actually the expression of the divergence of Ω in the basis $\{\Omega, \mathbf{u}, \mathbf{v}\}$). For the other two components of $\mathbf{r}_x(\Lambda)$, we perform exactly the same computations with a circular permutation of the roles of $\Omega, \mathbf{u}, \mathbf{v}$ to get $\mathbf{r}_x(\Lambda) \cdot \mathbf{u} = \nabla_x \cdot \mathbf{u}$ and $\mathbf{r}_x(\Lambda) \cdot \mathbf{v} = \nabla_x \cdot \mathbf{v}$. Therefore we obtain (4.28).

Finally we rewrite the equation for Λ as the evolution of the basis $\{\Omega, \mathbf{u}, \mathbf{v}\}$. To obtain the evolution of $\Lambda \mathbf{e}_k$ for $k = 1, 2, 3$, we multiply Eq. (4.22) by \mathbf{e}_k and compute to obtain:

$$\begin{aligned} \rho D_t \Omega + P_{\Omega^\perp} (c_3 \nabla_x \rho + c_4 \rho \mathbf{r}_x(\Lambda)) &= 0, \\ \rho D_t \mathbf{u} - \mathbf{u} \cdot (c_3 \nabla_x \rho + c_4 \rho \mathbf{r}_x(\Lambda)) \Omega + c_4 \rho \delta_x(\Lambda) \mathbf{v} &= 0, \\ \rho D_t \mathbf{v} - \mathbf{v} \cdot (c_3 \nabla_x \rho + c_4 \rho \mathbf{r}_x(\Lambda)) \Omega - c_4 \rho \delta_x(\Lambda) \mathbf{u} &= 0, \end{aligned}$$

where $D_t = \partial_t + c_2(\Omega \cdot \nabla_x)$. To perform the computations we have used for $\mathbf{w} = \nabla_x \rho$ or $\mathbf{w} = \mathbf{r}$ that

$$[\mathbf{w} \times \Omega]_\times \Omega = -P_{\Omega^\perp}(\mathbf{w}) \quad \text{and} \quad (\mathbf{w} \times \Omega) \times \mathbf{u} = (\mathbf{u} \cdot \mathbf{w}) \Omega,$$

since $\Omega \perp \mathbf{u}$ (analogously for \mathbf{v}). From here, using (4.28) we obtain straightforwardly Eqs. (4.29)–(4.31) for Ω, \mathbf{u} and \mathbf{v} , respectively. \square

5. Conclusions and Open Questions

In this work, we have presented a new flocking model through body attitude coordination. We have proposed an individual-based model where agents are described by their position and a rotation matrix (corresponding to the body attitude). From the individual-based model we have derived the macroscopic equations via the mean-field equations. We observe that the macroscopic equation gives rise to a new class of models, the SOHB. This model does not reduce to the more classical SOH, which is the continuum version of the Vicsek model. The dynamics of the SOHB system are more complex than those of the SOH ones of the Vicsek model. In a future work, we will carry out simulations of the individual-based model and the SOHB model and study the patterns that arise to compare them with the ones of the Vicsek and SOH model.

Also, there exist yet many open questions on the modelling side. For instance, one could consider that agents have a limited angle of vision, thus the so-called influence kernel K (see Sec. 3.1) is not isotropic any more, see Ref. 29 for the case of the Vicsek and SOH models. One could also consider a different interaction range for the influence kernel K that may give rise to a diffusive term in the macroscopic equations, see Ref. 19. Moreover, in the case of the SOH model, when the coordination frequency and noise intensity (quantities ν and D in the individual-based model (3.10)–(3.11)) are functions of the flux of the agents, then phase transitions occur at the macroscopic level¹⁹ (see also Refs. 4, 6, 20 and 44). An analogous feature is expected to happen in the present case. Finally, one could think of elaborating on the model by adding repulsive effects at short range and attraction effects at large range.

On the analytical side, this model opens also many questions like making Proposition 3.2 rigorous, which means dealing with stochastic differential equations with non-Lipschitz coefficients. In the context of the Vicsek model, the global well-posedness has been proven for the homogeneous mean-field Vicsek equation and also its convergence to the von Mises equilibria in Ref. 28, see also Ref. 31; an analogous result for our model will be desirable. The convergence of the Vicsek model to the model which was formally done in Ref. 24 has been recently achieved rigorously in Ref. 39. Again, one could also think of generalising these results to our case.

Appendix A. Special Orthogonal Group $\text{SO}(3)$

Throughout the text, we used repeatedly the following properties.

Proposition A.1. (Space decomposition in symmetric and antisymmetric matrices) *Denote by \mathcal{S} the set of symmetric matrices in \mathcal{M} and by \mathcal{A} the set of antisymmetric ones. Then*

$$\mathcal{S} \oplus \mathcal{A} = \mathcal{M} \quad \text{and} \quad \mathcal{A} \perp \mathcal{S}.$$

Proof. For $A \in \mathcal{M}$ we have $A = \frac{1}{2}(A + A^T) + \frac{1}{2}(A - A^T)$, the first term being symmetric and the second antisymmetric. The orthogonality comes from the properties of the trace, namely $\text{tr}(A^T) = \text{tr}(A)$, and $\text{tr}(AB) = \text{tr}(BA)$ for $B \in \mathcal{M}$. Indeed if $P \in \mathcal{A}$ and $S \in \mathcal{S}$ then $\text{tr}(P^T S) = \text{tr}(SP^T) = \text{tr}(PS^T) = -\text{tr}(P^T S)$. Hence $P \cdot S = \frac{1}{2}\text{tr}(P^T S) = 0$. \square

Proposition A.2. (Tangent space to $\text{SO}(3)$) *For $A \in \text{SO}(3)$, denote by T_A the tangent space to $\text{SO}(3)$ at A . Then*

$$M \in T_A \quad \text{if and only if there exists } P \in \mathcal{A} \text{ such that } M = AP,$$

or equivalently the same statement with $M = PA$. Consequently, we have that

$$M \in T_A^\perp \quad \text{if and only if there exists } S \in \mathcal{S} \text{ such that } M = AS,$$

or equivalently the same statement with $M = SA$.

Proof. We have that $M \in T_A$ if and only if there exists a curve $\Lambda(t)$ from the neighbourhood of 0 in \mathbb{R} to $\text{SO}(3)$ such that $\Lambda(0) = A$ and $\Lambda'(0) = M$. We then have

$$\begin{aligned} \text{Id} &= \Lambda(t)\Lambda^T(t) = (A + tM + o(t))(A^T + tM^T + o(t)) \\ &= \text{Id} + t(A^T M + M^T A) + o(t). \end{aligned}$$

So if $M \in T_A$, we must have $(A^T M + M^T A) = 0$, that is to say that $P = A^T M \in \mathcal{A}$.

Conversely if $M = AP$ with $P \in \mathcal{A}$, the solution of the linear differential equation $\Lambda'(t) = \Lambda(t)P$ with $\Lambda(0) = A$ is given by $\Lambda(t) = Ae^{tP}$ so it is a curve

in $\text{SO}(3)$. Indeed we have $\Lambda(t)^T \Lambda = (e^{tP})^T e^{tP} = e^{tP^T} e^{tP} = e^{-tP} e^{tP} = \text{Id}$. Since $\Lambda'(0) = AP = M$, we get that $M \in T_A$. The equivalent condition comes from the fact that if $M = AP$, with $P \in \mathcal{A}$, then $M = APA^T A = \tilde{P}A$ with $\tilde{P} \in \mathcal{A}$. Finally the last part is obtained thanks to Proposition A.1) and the fact that the dot product is left (and right) invariant with respect to $\text{SO}(3)$: if $B, C \in \mathcal{M}$ and $A \in \text{SO}(3)$, then $AB \cdot AC = \frac{1}{2} \text{tr}(B^T A^T AC) = B \cdot C$. \square

Proposition A.3. (Projection operator on the tangent space) *Let $A \in \text{SO}(3)$ and $M \in \mathcal{M}$ (set of square matrices). Let P_{T_A} be the orthogonal projection on T_A (tangent space at A), then*

$$P_{T_A}(M) = \frac{1}{2}(M - AM^T A). \quad (\text{A.1})$$

Notice that then

$$P_{T_A^\perp}(M) = \frac{1}{2}(M + AM^T A).$$

Proof. It suffices to verify that the expression given for $P_{T_A}(M)$ satisfies $P_{T_A}(M) \in T_A$ and $M - P_{T_A}(M) \in T_A^\perp$, that is to say $A^T P_{T_A} M \in \mathcal{A}$ and $A^T(M - P_{T_A}(M)) \in \mathcal{S}$ thanks to Proposition A.2. We have indeed $A^T \frac{1}{2}(M - AM^T A) = \frac{1}{2}(A^T M - M^T A)$ which is clearly antisymmetric, and $A^T \frac{1}{2}(M + AM^T A) = \frac{1}{2}(A^T M + M^T A)$ which is symmetric. \square

To compute the gradient in $\text{SO}(3)$ of a function $\psi : \text{SO}(3) \rightarrow \mathbb{R}$ we will consider $A(\varepsilon)$ a differentiable curve in $\text{SO}(3)$ such that

$$A(0) = A, \quad \left. \frac{d}{d\varepsilon} A(\varepsilon) \right|_{\varepsilon=0} = \delta_A \in T_A,$$

then $\nabla_A \psi(A)$ is the element of T_A such that for any $\delta_A \in T_A$, we have

$$\lim_{\varepsilon \rightarrow 0} \frac{\psi(A(\varepsilon)) - \psi(A)}{\varepsilon} = \nabla_A \psi(A) \cdot \delta_A.$$

In particular, one can check that

$$\nabla_A(A \cdot M) = P_{T_A}(M), \quad M \in \mathcal{M}. \quad (\text{A.2})$$

We now show that the differential equation given by this gradient has trajectories supported on geodesics.

Proposition A.4. *If $B \in \text{SO}(3)$ and $A_0 \in \text{SO}(3)$, the trajectory of the solution of the differential equation $\frac{dA}{dt} = \nu(A \cdot B) P_{T_A} B = \nu(A \cdot B) \nabla_A(A \cdot B)$ with $A(0) = A_0$ (and with ν smooth and positive) is supported on a geodesic from A_0 to B .*

Proof. Indeed, write $B^T A_0 = \exp(\theta_0 [\mathbf{n}]_\times)$ thanks to Rodrigues' formula (3.4) with $[\mathbf{n}]_\times$ an antisymmetric matrix of unit norm and $\theta_0 \in [0, \pi]$. If we set

$A(t) = B \exp(\theta(t)[\mathbf{n}]_{\times})$ where θ satisfies the equation $\theta' = -\nu(\frac{1}{2} + \cos \theta) \sin \theta$ with $\theta(0) = \theta_0$, we get

$$\begin{aligned} \frac{dA}{dt} &= B \exp(\theta(t)[\mathbf{n}]_{\times}) \theta'(t) [\mathbf{n}]_{\times} \\ &= -\nu \left(\frac{1}{2} + \cos \theta(t) \right) B \exp(\theta(t)[\mathbf{n}]_{\times}) \sin \theta(t) [\mathbf{n}]_{\times}. \end{aligned}$$

Now, thanks to the expression (3.3), we have

$$\sin \theta [\mathbf{n}]_{\times} = \frac{1}{2} (\exp(\theta [\mathbf{n}]_{\times}) - \exp(\theta [\mathbf{n}]_{\times})^T) = \frac{1}{2} (B^T A - A^T B),$$

and $A \cdot B = \text{Id} \cdot AB^T = \frac{1}{2} \text{tr}(\exp(\theta [\mathbf{n}]_{\times})) = \frac{1}{2} + \cos \theta$ thanks to (3.1). Therefore we obtain

$$\frac{dA}{dt} = -\nu(A \cdot B) A \frac{1}{2} (B^T A - A^T B) = \nu(A \cdot B) P_{T_A} B,$$

thanks to (A.1) and we have $A(0) = A_0$. Since $\theta \in [0, \theta_0] \mapsto \exp(\theta [\mathbf{n}]_{\times})$ is a geodesic between Id and $B^T A_0$, then $\theta \mapsto B \exp(\theta [\mathbf{n}]_{\times})$ is a geodesic between B and A , and the solution $A(t)$ is supported on this geodesic. It is also easy to see that, except in the case $\theta_0 = \pi$ or $\theta_0 = 0$, for which the solution is constant, the function $t \mapsto \theta(t)$ (solution of the one-dimensional differential equation

$$\theta' = -\nu \left(\frac{1}{2} + \cos \theta \right) \sin \theta$$

is positive, decreasing, and converge exponentially fast to 0, with an asymptotic exponential rate $\nu(\frac{3}{2})$. Therefore, as time goes to infinity, the trajectory covers the whole geodesic from A_0 to B (excluded). \square

We now turn to the proofs of the expressions of the gradient, the volume form and the divergence in $\text{SO}(3)$ in the so-called Euler axis-angle coordinates, that were presented in Sec. 4.2.

Proof of Proposition 4.1. Consider a curve in $\text{SO}(3)$ given by

$$A(t) = \exp(\theta(t)[\mathbf{n}]_{\times}(t)) = \text{Id} + \sin(\theta(t))[\mathbf{n}]_{\times}(t) + (1 - \cos(\theta(t)))[\mathbf{n}]_{\times}^2(t)$$

(following (3.3)–(3.2)) with $A(0) = A$, $\theta(0) = \theta$ and $[\mathbf{n}]_{\times}(t) = [\mathbf{n}(t)]_{\times}$, $\mathbf{n}(0) = \mathbf{n}$. Define

$$\delta_A = A'(0) \in T_A,$$

$$\delta_{\theta} = \theta'(0) \in \mathbb{R},$$

$$\delta_{\mathbf{n}} = \mathbf{n}'(0),$$

$$\delta_{[\mathbf{n}]_{\times}} = [\mathbf{n}]'_{\times}(0) = [\delta_{\mathbf{n}}]_{\times}.$$

With these notations, for a function $f = f(A(\theta, \mathbf{n}))$ it holds

$$\nabla_A f \cdot \delta_A = \frac{\partial f}{\partial \theta} \delta_\theta + \nabla_{\mathbf{n}} f \cdot \delta_{\mathbf{n}}. \quad (\text{A.3})$$

On the other hand, it holds true that

$$\begin{aligned} \delta_A &= A[\mathbf{n}]_\times \delta_\theta + \sin \theta \delta_{[\mathbf{n}]_\times} + (1 - \cos \theta) ([\mathbf{n}]_\times \delta_{[\mathbf{n}]_\times} + \delta_{[\mathbf{n}]_\times} [\mathbf{n}]_\times) \\ &= A[\mathbf{n}]_\times \delta_\theta + A A^T (\sin \theta \delta_{[\mathbf{n}]_\times} + (1 - \cos \theta) ([\mathbf{n}]_\times \delta_{[\mathbf{n}]_\times} + \delta_{[\mathbf{n}]_\times} [\mathbf{n}]_\times)) \\ &= A[\mathbf{n}]_\times \delta_\theta + A (\text{Id} - \sin \theta [\mathbf{n}]_\times + (1 - \cos \theta) [\mathbf{n}]_\times^2) \\ &\quad \times (\sin \theta \delta_{[\mathbf{n}]_\times} + (1 - \cos \theta) ([\mathbf{n}]_\times \delta_{[\mathbf{n}]_\times} + \delta_{[\mathbf{n}]_\times} [\mathbf{n}]_\times)) \\ &= A[\mathbf{n}]_\times \delta_\theta + A (\sin \theta \delta_{[\mathbf{n}]_\times} + (1 - \cos \theta) (\delta_{[\mathbf{n}]_\times} [\mathbf{n}]_\times - [\mathbf{n}]_\times \delta_{[\mathbf{n}]_\times})) \\ &= A[\mathbf{n}]_\times \delta_\theta + 2 \sin(\theta/2) A (\cos(\theta/2) [\delta_{\mathbf{n}}]_\times + \sin(\theta/2) [\mathbf{n} \times \delta_{\mathbf{n}}]_\times) \\ &= A[\mathbf{n}]_\times \delta_\theta + L_{[\mathbf{n}]_\times} (\delta_{[\mathbf{n}]_\times}), \end{aligned} \quad (\text{A.4})$$

where the last line defines $L_{[\mathbf{n}]_\times}$. In the first line, the term in δ_θ is obtained by differentiating the exponential form (3.4) of $A(t)$ assuming that $[\mathbf{n}]_\times(t)$ is constant. The term in $\delta_{[\mathbf{n}]_\times}$ is obtained by differentiating Rodrigues' formula (3.3). To do the computation we have used Rodrigues' formula (3.3) to express A^T and the facts that $[\mathbf{n}]_\times^3 = -[\mathbf{n}]_\times$; $[\mathbf{n}]_\times \delta_{[\mathbf{n}]_\times} [\mathbf{n}]_\times = 0$; and $\delta_{[\mathbf{n}]_\times} [\mathbf{n}]_\times - [\mathbf{n}]_\times \delta_{[\mathbf{n}]_\times} = [\delta_{\mathbf{n}} \times \mathbf{n}]_\times$.

In particular notice that $\{[\mathbf{n}]_\times, [\delta_{\mathbf{n}}]_\times, [\mathbf{n} \times \delta_{\mathbf{n}}]_\times\}$ is an orthogonal basis of \mathcal{A} (antisymmetric matrices) from which we obtain a basis of T_A (by Proposition A.2). So, we just need to compute the components of $\nabla_A f$ in $\text{span}\{A[\mathbf{n}]_\times\}$ and $\text{span}\{(A[\mathbf{n}]_\times)^\perp\}$.

We will show that the component in $\text{span}\{A[\mathbf{n}]_\times\}$ is given by

$$P_{A[\mathbf{n}]_\times}(\nabla_A f) = \frac{\partial f}{\partial \theta} A[\mathbf{n}]_\times \quad (\text{A.5})$$

and the one on $\text{span}\{(A[\mathbf{n}]_\times)^\perp\}$ is

$$P_{(A[\mathbf{n}]_\times)^\perp}(\nabla_A f) = \frac{1}{2 \sin(\theta/2)} A (\cos(\theta/2) [\nabla_{\mathbf{n}} f]_\times + \sin(\theta/2) [\mathbf{n} \times \nabla_{\mathbf{n}} f]_\times). \quad (\text{A.6})$$

The sum of the two previous expressions gives (4.2):

$$\nabla_A f = P_{A[\mathbf{n}]_\times}(\nabla_A f) + P_{(A[\mathbf{n}]_\times)^\perp}(\nabla_A f).$$

The component (A.5) is computed considering the case where $\delta_{\mathbf{n}} = 0$ in (A.4)–(A.3), so that

$$\nabla_A f \cdot \delta_A = \nabla_A f \cdot A[\mathbf{n}]_\times \delta_\theta = \frac{\partial f}{\partial \theta} \delta_\theta.$$

Expression (A.5) is obtained by noticing that

$$(A[\mathbf{n}]_\times) \cdot (A[\mathbf{n}]_\times) = [\mathbf{n}]_\times \cdot [\mathbf{n}]_\times = \mathbf{n} \cdot \mathbf{n} = 1 \quad (\text{using (3.7)}).$$

To obtain the component (A.6), consider the case $\delta_\theta = 0$ in (A.4) and (A.3) so that

$$\nabla_A f \cdot \delta_A = \nabla_A f \cdot L_{[\mathbf{n}]_\times}(\delta_{[\mathbf{n}]_\times}) = \nabla_{\mathbf{n}} f \cdot \delta_{\mathbf{n}}, \quad (\text{A.7})$$

where $L_{[\mathbf{n}]_\times}$ is given in (A.4).

We have that

$$P_{(A[\mathbf{n}]_\times)^\perp}(\nabla_A f) = A[\mathbf{u}]_\times \quad \text{for some } \mathbf{u} \perp \mathbf{n}.$$

The goal is to compute \mathbf{u} as a function of $\mathbf{v} := \nabla_{\mathbf{n}} f$. By (A.7) we have that

$$A[\mathbf{u}]_\times \cdot L_{[\mathbf{n}]_\times}(\delta_{[\mathbf{n}]_\times}) = \nabla_{\mathbf{n}} f \cdot \delta_{\mathbf{n}}.$$

This implies that

$$2 \sin(\theta/2)[\mathbf{u}]_\times \cdot (\cos(\theta/2)[\delta_{\mathbf{n}}]_\times + \sin(\theta/2)[\mathbf{n} \times \delta_{\mathbf{n}}]_\times) = \mathbf{v} \cdot \delta_{\mathbf{n}} \quad \text{for all } \delta_{\mathbf{n}} \perp \mathbf{n},$$

so (see (3.7)) we get

$$2 \sin(\theta/2)(\cos(\theta/2)\mathbf{u} + \sin(\theta/2)\mathbf{u} \times \mathbf{n}) \cdot \delta_{\mathbf{n}} = \mathbf{v} \cdot \delta_{\mathbf{n}}.$$

Since this is true for all $\delta_{\mathbf{n}}$ orthogonal to \mathbf{n} , we get

$$\mathbf{v} = 2 \sin(\theta/2)(\cos(\theta/2)\mathbf{u} + \sin(\theta/2)\mathbf{u} \times \mathbf{n}).$$

From here we can get the expression of $\mathbf{n} \times \mathbf{v}$ in terms of \mathbf{u} and $\mathbf{n} \times \mathbf{u}$. After some computations we finally obtain that

$$\mathbf{u} = \frac{1}{2 \sin(\theta/2)}(\cos(\theta/2)\mathbf{v} + \sin(\theta/2)\mathbf{n} \times \mathbf{v}). \quad \square$$

Proof of the volume form, Lemma 4.2. We denote by g the metric of the Riemannian manifold $\text{SO}(3)$ associated to the inner product

$$A \cdot B = \frac{1}{2} \text{tr}(A^T B), \quad A, B \in \text{SO}(3).$$

The volume form is proportional to $\sqrt{\det(g)}$.³⁰ We compute the volume form using spherical coordinates, i.e. we consider the coordinates

$$(\theta, \phi, \psi) \in [0, \pi] \times [0, \pi] \times [0, 2\pi].$$

Given the Euler axis-angle coordinates (θ, \mathbf{n}) we have that

$$\mathbf{n} = \begin{pmatrix} \sin \phi \cos \psi \\ \sin \phi \sin \psi \\ \cos \phi \end{pmatrix}.$$

For the spherical coordinate system, we consider the vector field $(\frac{\partial}{\partial \theta}, \frac{\partial}{\partial \phi}, \frac{\partial}{\partial \psi})$. Denoting

$$Y_1 = \frac{\partial A}{\partial \theta}, \quad Y_2 = \frac{\partial A}{\partial \phi}, \quad Y_3 = \frac{\partial A}{\partial \psi}, \quad A \in \text{SO}(3),$$

we get that $(Y_i)_{i=1,2,3} \in T_A(\text{SO}(3))$ forms a basis of vectors fields at A .

The metric g is defined as $g_{ij} = g(Y_i, Y_j) = \frac{1}{2} \text{tr}(Y_i^T Y_j)$, $i, j = 1, 2, 3$. We compute next each term. First, we know that for a given $\delta_A \in T_A$, there exist $\delta_\theta, \delta_\psi, \delta_\phi \in \mathbb{R}$

such that

$$\delta_A = \frac{\partial A}{\partial \theta} \delta_\theta + \frac{\partial A}{\partial \phi} \delta_\phi + \frac{\partial A}{\partial \psi} \delta_\psi$$

and also for a given $\delta_{\mathbf{n}} \in T_{\mathbf{n}}(S^2)$ (the tangent plane to the sphere at \mathbf{n}), there exist $\delta'_\psi, \delta'_\phi$ such that

$$\delta_{\mathbf{n}} = \frac{\partial \mathbf{n}}{\partial \phi} \delta'_\phi + \frac{\partial \mathbf{n}}{\partial \psi} \delta'_\psi.$$

Now, following the computation given in (A.4) we have that for $\delta_\theta = 1, \delta_\phi = 0, \delta_\psi = 0$:

$$\frac{\partial A}{\partial \theta} = \delta_A = A[\mathbf{n}]_\times.$$

Now, if $\delta_\theta = 0, \delta_\phi = 1, \delta_\psi = 0$ then, using that $\delta_{\mathbf{n}} = \frac{\partial \mathbf{n}}{\partial \phi}$ we have that

$$\frac{\partial A}{\partial \phi} = \delta_A = 2 \sin(\theta/2) A \left[R_{\mathbf{n}, \theta/2} \left(\frac{\partial \mathbf{n}}{\partial \phi} \right) \right]_\times,$$

where

$$R_{\mathbf{n}, \theta/2}(\mathbf{v}) = \cos(\theta/2) \mathbf{v} + \sin(\theta/2) (\mathbf{n} \times \mathbf{v}),$$

which corresponds to the rotation of the vector \mathbf{v} around \mathbf{n} by an angle $\theta/2$ (anti-clockwise) as long as $\mathbf{v} \cdot \mathbf{n} = 0$. Analogously one can also deduce that

$$\frac{\partial A}{\partial \psi} = 2 \sin(\theta/2) \left[R_{\mathbf{n}, \theta/2} \left(\frac{\partial \mathbf{n}}{\partial \psi} \right) \right]_\times.$$

From here, using that $\|\frac{\partial \mathbf{n}}{\partial \phi}\|^2 = 1$ and $\|\frac{\partial \mathbf{n}}{\partial \psi}\|^2 = \sin^2 \phi$, we conclude that

$$g = \begin{pmatrix} 1 & 0 & 0 \\ 0 & 4 \sin^2(\theta/2) & 0 \\ 0 & 0 & 4 \sin^2(\theta/2) \sin^2 \phi \end{pmatrix}.$$

Notice that to compute $g(\frac{\partial A}{\partial \theta}, \frac{\partial A}{\partial \phi})$ we use that $R_{\mathbf{n}, \theta/2}(\frac{\partial \mathbf{n}}{\partial \phi}) \perp \mathbf{n}$.

Finally we have that

$$\sqrt{\det(g)} = 4 \sin^2(\theta/2) \sin \phi$$

and therefore

$$\begin{aligned} \int_{\text{SO}(3)} f(A) dA &= \int_{[0, \pi] \times [0, \pi] \times [0, 2\pi]} \tilde{f}(\theta, \phi, \psi) 4 \sin^2(\theta/2) \sin \phi d\theta d\phi d\psi \\ &= 4 \int_{\theta \in [0, \pi]} \left(\int_{[0, \pi] \times [0, 2\pi]} \tilde{f}(\theta, \phi, \psi) \sin \phi d\phi d\psi \right) \sin^2(\theta/2) d\theta. \end{aligned}$$

The term $\sin \phi d\phi d\psi$ is the volume element in the sphere S^2 so we have that

$$\int_{S^2} \hat{f}(\theta, \mathbf{n}) d\mathbf{n} = \int_{[0, \pi] \times [0, 2\pi]} \tilde{f}(\theta, \phi, \psi) \sin \phi d\phi d\psi.$$

Therefore, the volume element corresponding to the Euler axis-angle coordinates is proportional to $\sin^2(\theta/2) d\theta d\mathbf{n}$. Since the volume element is defined up to a constant, we choose the constant c such that

$$\int_0^{2\pi} c \sin^2(\theta/2) d\theta = 1,$$

i.e. $c = 2/\pi$. In conclusion, the volume element in the Euler axis-angle coordinates corresponds to

$$\frac{2}{\pi} \sin^2(\theta/2) d\theta d\mathbf{n}.$$

□

Proof of divergence formula, Proposition 4.2. We compute the divergence by duality of the gradient, Proposition 4.1. Let $f = f(A)$ be a function and consider

$$\begin{aligned} & \int_{\text{SO}(3)} \nabla_A \cdot B(A) f(A) dA \\ &= - \int_{\text{SO}(3)} B(A) \cdot \nabla_A f(A) dA \\ &= - \int_{(0, \pi) \times S^2} W(\theta) b(\theta, \mathbf{n}) \partial_\theta f(\theta, \mathbf{n}) d\theta d\mathbf{n} \\ &\quad - \int_{(0, \pi) \times S^2} \frac{W(\theta)}{2 \sin(\theta/2)} \mathbf{v}(\theta, \mathbf{n}) \cdot (\cos(\theta/2) \nabla_{\mathbf{n}} f(\mathbf{n}, \theta) + \sin(\theta/2) \mathbf{n} \\ &\quad \times \nabla_{\mathbf{n}} f(\mathbf{n}, \theta)) d\theta d\mathbf{n} \\ &= \int_{(0, \pi) \times S^2} \frac{f(\theta, \mathbf{n})}{\sin^2(\theta/2)} \partial_\theta (\sin^2(\theta/2) b(\theta, \mathbf{n})) W(\theta) d\theta d\mathbf{n} \\ &\quad + \int_{(0, \pi) \times S^2} \frac{f(\theta, \mathbf{n})}{2 \sin(\theta/2)} \nabla_{\mathbf{n}} \cdot (\mathbf{v}(\theta, \mathbf{n}) \cos(\theta/2) + \sin(\theta/2) (\mathbf{v}(\theta, \mathbf{n}) \times \mathbf{n})) \\ &\quad \times W(\theta) d\theta d\mathbf{n}, \end{aligned}$$

where W is given by (4.3), from which we deduce the result. □

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Data Statement

No new data was generated in the course of this research.

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